

FREE VIBRATIONS IN A WAVE EQUATION MODELING MEMS*

CARLOS GARCÍA-AZPEITIA[†] AND JEAN-PHILIPPE LESSARD[‡]

Abstract. We study a nonlinear wave equation appearing as a model for a membrane (without viscous effects) under the presence of an electrostatic potential with strength λ . The membrane has a unique stable branch of steady states u_λ for $\lambda \in [0, \lambda_*]$. We prove that the branch u_λ has an infinite number of branches of periodic solutions (free vibrations) bifurcating when the parameter λ is varied. Furthermore, using a functional setting, we compute numerically the branch u_λ and their branches of periodic solutions. This approach is useful to validate rigorously the steady states u_λ at the critical value λ_* .

Dedicated to the memory of G. Flores.

Key words. microelectromechanical system, wave equation, periodic solutions

AMS subject classifications. 35B10 35B32 35L81

1. Introduction. We consider an idealized device that consists of an elastic plate suspended above a rigid ground plate. This device falls in the category of microelectromechanical systems (MEMS). The membrane is taken to be rectangular with two fixed parallel sides, while the other sides are considered to be thin and free. When a potential difference is applied between the membrane and the plate, the membrane deflects towards the ground plate. We assume that dissipation which might result from viscous effects on the moving membrane can be neglected. Under these assumptions, the deformation of the elastic membrane is described by the dimensionless equation

$$(1.1) \quad U_{tt} - U_{xx} + \frac{\lambda}{(1+U)^2} = 0, \quad x \in [-\pi/2, \pi/2],$$

where $U(x, t)$ satisfies the Dirichlet boundary conditions $U(\pm\pi/2) = 0$, and the parameter λ represents the strength of the applied voltage. A derivation of the nonlinearity leading to a general equation modeling the electrostatic membrane

$$\varepsilon^2 U_{tt} + \nu U_t - U_{xx} + \frac{\lambda f(x)}{(1+U)^2} = 0,$$

for which the nonlinear wave equation (1.1) is a particular case, can be found in [27], where the parameter ε represents the strength of the inertial term, ν the viscosity and the function $f(x)$ encodes the dielectric permittivity of the membrane.

In the design of microelectronic devices, it is relevant to study whether or not the membrane touches the ground plate. This phenomenon is called *touchdown* or *quenching*. Mathematically, quenching occurs if there is a point (x, t) such that $U(x, t) = -1$. A vast literature exists on the study of MEMS via parabolic and hyperbolic PDE modelling. Equation (1.1) is in fact a special case of the more general MEMS parabolic and hyperbolic PDE models, and a vast mathematical literature is dedicated to their study. Let us give a few examples.

*The second author was funded by NSERC.

[†] Depto. Matemáticas y Mecánica IIMAS, Universidad Nacional Autónoma de México, Apdo. Postal 20-726, 01000 Ciudad de México, México, (cgazpe@mym.iimas.unam.mx).

[‡] McGill University, Department of Mathematics and Statistics, 805 Sherbrooke Street West, Montreal, QC, H3A 0B9, Canada (jp.lessard@mcgill.ca), (<http://www.math.mcgill.ca/jplessard/>).

The local existence of solutions and the existence of quenching at a finite time for a parabolic equation modelling MEMS is analyzed in [12, 18, 19], and references therein. The case of a nonlocal parabolic equation modelling MEMS is proposed in [13]. The existence of solutions and the finite-time quenching for a damped wave equation modelling MEMS is analyzed in [11]. Generalizations of the wave equation (1.1) have been studied in [5, 22, 31].

The study of periodic orbits in MEMS models have also received their fair share of attention. Their relevance comes from the fact that they persist as small oscillations with no quenching. In [21], periodic solutions were observed numerically by solving an initial value problem for a non-local wave equation modelling MEMS. The study of periodic solutions in Hamiltonian PDEs (such as equation (1.1)) presents intrinsic problems associated to infinite-dimensional kernels [28], lack of compactness [10], or small divisor problems [7]. The small divisor problem was avoided in [1, 24, 28] by imposing restrictions on the temporal period of the solutions of a nonlinear wave equation. By imposing similar restrictions, the articles [17] obtains the existence of continuous branches of periodic solutions for a nonlinear wave equation in a sphere, and in [6, 15] for a Hamiltonian PDE appearing in the n -vortex filament problem.

The equation (1.1) has a family of stable steady states u_λ for $\lambda \in [0, \lambda_*]$. For instance, the existence of steady states of (1.1) in the N -dimension ball was proven in [16]. In the present work, we prove the existence of continuous families of periodic solutions near the branch of steady states (Theorem 1.1). This is our main contribution, and to the best of our knowledge, this result is new. The main challenge encountered when proving the existence of the periodic solutions is that the trivial branch u_λ and its associated spectrum are not known explicitly. Indeed, while the existence of continuous families of periodic solutions has been obtained before for Hamiltonian PDEs in [6, 15, 17, 24], in those articles the trivial branch and the spectrum of its linearized equation are known explicitly. We overcome this problem with delicate estimates of the spectrum which depend on an accurate estimate for the steady state u_{λ_*} at the critical value λ_* . While the critical value λ_* is known (see for instance [13, 22]), no estimate for the steady state u_{λ_*} is known. That leads to our second main contribution, which is to obtain precise and rigorous estimate on the steady state u_{λ_*} at the critical value λ_* (Theorem 1.2). Finally, our third contribution is to present a systematic approach to compute numerically the families of periodic solutions using Chebyshev series expansion (in space) and Fourier series expansions (in time) (see Figures 2, 3, 4 and 5).

Specifically, the linear operator of the stationary equation (1.1) at u_λ is given by

$$(1.2) \quad A(\lambda)u \stackrel{\text{def}}{=} -\partial_x^2 u - \frac{2\lambda}{(1+u_\lambda)^3} u : H_0^2 \subset L^2 \rightarrow L^2,$$

where $H_0^2([-\pi/2, \pi/2]; \mathbb{R})$ is the Sobolev space of functions $u(x)$ satisfying Dirichlet boundary conditions $u(\pm\pi/2) = 0$. The operator A is self-adjoint and positive definite for $\lambda \in [0, \lambda_*)$ because the first eigenvalue $\mu_1(\lambda)$ is positive (Theorem 4.2 in [16]). Thus A has eigenvalues $0 < \mu_1(\lambda) < \mu_2(\lambda) < \dots$ with eigenfunctions satisfying

$$(1.3) \quad A(\lambda)v_k(x; \lambda) = \mu_k(\lambda)v_k(x; \lambda), \quad k \in \mathbb{N}.$$

Our main theorem regarding the existence of periodic solutions is the following.

THEOREM 1.1. *There is an infinite number of non-resonant parameters $\lambda_0 \in (0, \lambda_*)$, associated with numbers $p, q, k \in \mathbb{N}$ by the relation $\mu_k(\lambda_0) = (p/q)^2$, such*

that there is a local continuum of $2\pi q/p$ -periodic solutions bifurcating from the steady state $u_\lambda(x)$ with $\lambda = \lambda_0$. The local bifurcation consists of free vibrations satisfying the estimates

$$(1.4) \quad \begin{aligned} U(t, x; \lambda) &= u_{\lambda_0}(x) + b \cos(pt/q) v_k(x; \lambda_0) + \mathcal{O}_{C^4}(b^2), \\ \lambda &= \lambda_0 + \mathcal{O}(b), \end{aligned}$$

where $b \in [0, b_0]$ represents a parametrization of the local bifurcation for some $b_0 > 0$ and $\mathcal{O}_{C^4}(b^2)$ is a function of order b^2 in the C^4 -norm. Furthermore, the bifurcation has symmetries

$$\begin{aligned} U(t, x) &= U(-t, x) = U(t, -x), & \text{for } k \text{ odd,} \\ U(t, x) &= U(-t, x) = U(t + p\pi/q, -x), & \text{for } k \text{ even.} \end{aligned}$$

The existence of a branch of periodic solutions arising from the family of steady states u_λ is set as a branch of zeros for the functional equation

$$(1.5) \quad L(\lambda)u + g(u; \lambda) = 0, \quad L(\lambda) \stackrel{\text{def}}{=} (p/q)^2 \partial_t^2 + A(\lambda),$$

where $A(\lambda)$ is given in (1.2) and $g = \mathcal{O}(u^2)$ is an analytic nonlinearity defined in a neighborhood of zero. Here $u \in H_{sym}^s$ represents a perturbation from the steady state u_λ , where $H_{sym}^s(S^1 \times [-\pi/2, \pi/2]; \mathbb{R})$ is the Sobolev space of even 2π -periodic functions $u(t, x)$ satisfying Dirichlet boundary conditions $u(t, \pm\pi/2) = 0$. The spectrum of the elliptic operator $A(\lambda)$ is not explicit, but it can be estimated by applying the Courant-Fischer-Weyl minmax theorem (Chapter 39 in [29]). These estimates are essential to show that the linear operator $L(\lambda_0)$ has a finite-dimensional kernel, which is non-trivial for a dense set of values λ_0 in the interval $(0, \lambda_*)$. We implement a Lyapunov-Schmidt reduction for equation (1.5). The range equation is solved by the contracting mapping theorem and by proving that the linear operator is invertible (but not compact) in the range. The bifurcation equation is solved for a non-resonant value $\lambda_0 \in (0, \lambda_*)$ using the Crandall-Rabinowitz theorem [8]. The proof of the main theorem is finished by showing that the number of non-resonant points λ_0 in $(0, \lambda_*)$ is infinite. It is important to mention that the main theorem in [25] and Remark 2.17 imply that the set of bifurcation points λ_0 is dense in $(0, \lambda_*)$, but the bifurcations arising from these possibly resonant bifurcation points do not satisfy the estimates or symmetries of our main theorem.

The proof of Theorem 1.1 requires estimating rigorously the minimum value of the steady state u_λ at the critical value $\lambda = \lambda_*$. The proof of a precise estimate for the steady state u_{λ_*} is computer-assisted and is done independently using a Newton-Kantorovich argument based on the radii polynomial approach (e.g. see [9, 20, 32] and the references therein). To obtain a rigorous control on u_{λ_*} , we use Chebyshev polynomials series expansions. Since the Chebyshev polynomials are naturally defined on the interval $[-1, 1]$, we rescale the space domain $[-\pi/2, \pi/2]$ to $[-1, 1]$. Specifically, let $\tilde{U}(t, x; \lambda)$ be a solution of equation (1.1), then the scaled function $U(t, y; \lambda) \stackrel{\text{def}}{=} \tilde{U}(\alpha t, x; \alpha^{-2}\lambda)$ with $\alpha = \pi/2$ and $x = \alpha y$, is a solution of the equation

$$(1.6) \quad U_{tt} - U_{yy} + \frac{\lambda}{(1+U)^2} = 0, \quad y \in [-1, 1], \quad U(\pm 1) = 0.$$

We denote the critical value of the scaled equation (1.6) by $\lambda^* \stackrel{\text{def}}{=} (\pi/2)^{-2} \lambda_*$. The articles [13, 22] study the model in the domain $[-1/2, 1/2]$ and get an exact implicit formula for the critical value $4\lambda^*$, which later is approximated numerically by the value 1.400016469. Using the radii polynomial approach, we prove the following result.

THEOREM 1.2. *The branch of steady states u_λ of the MEMS equation (1.6) undergoes a saddle-node bifurcation at the steady state u_{λ^*} satisfying the precise estimates*

$$|u_{\lambda^*} - \bar{u}_{\lambda^*}|_{C^0[-1,1]} \leq r = 5.7 \times 10^{-12},$$

where \bar{u}_{λ^*} is the numerical approximation portrayed in Figure 1, and such that

$$\begin{aligned} \lambda^* &\in 0.350004119342744 + [-r, r], \\ u_{\lambda^*}(0) &\in -0.388346718912783 + [-r, r]. \end{aligned}$$

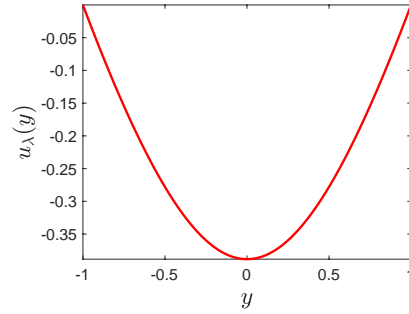


FIG. 1. *The numerical approximation $\bar{u}_{\lambda^*}(t)$ obtained in Section 3 for the steady state of equation (1.6) at the critical value λ^* .*

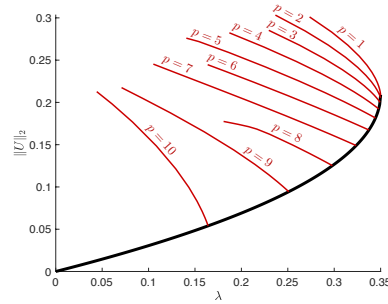


FIG. 2. *Several branches of periodic solutions for $k = 1$, $q = 11$ and $p \in \{1, \dots, 10\}$.*

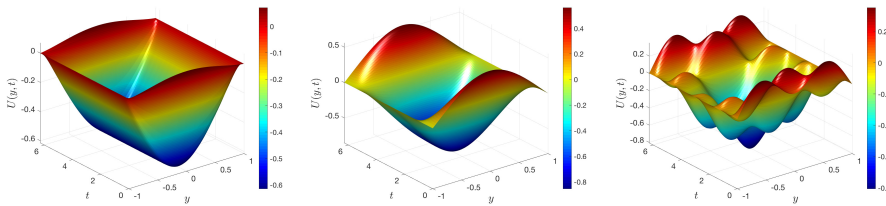


FIG. 3. *Periodic solutions on the branches appearing on Figure 2 for $p = 1$ at $\lambda = 0.273478006926454$ (left), $p = 7$ at $\lambda = 0.105151105978289$ (center), and $p = 8$ at $\lambda = 0.180819225410784$ (right).*

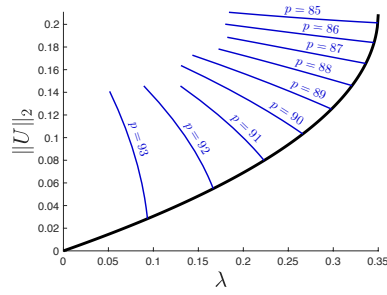


FIG. 4. Several branches of periodic orbits for $k = 2$, $q = 47$ and $p \in \{85, \dots, 93\}$.

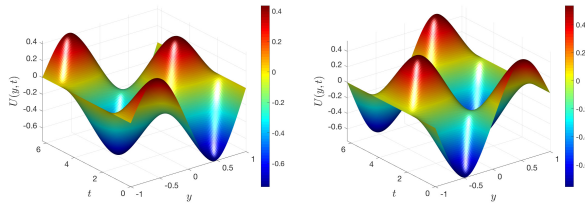


FIG. 5. The periodic solutions on the branches $p = 88$ at $\lambda = 0.172417160845949$ (left), and $p = 90$ at $\lambda = 0.130674749789634$ (right) of Figure 4. Here, $k = 2$ and $q = 47$.

The proof of this theorem is performed by solving the boundary value problem (4.6) with polynomial nonlinearities. This setting simplifies the estimates of the computer-assisted proof by considering the Banach algebra property of spaces of Chebyshev sequences with geometric decay (e.g. see [20, 32]). Note that the computations of the branch of steady states and the spectrum of the linear elliptic operator are also performed using Chebyshev series and are set as polynomial boundary value problems given in (3.2) and (3.9), respectively. For the periodic solutions we use the fact that $U(t, x)$ is a $2\pi q/p$ -periodic solution of equation (1.1) if and only if the rescaled function $U(t, y)$ is a $4q/p$ -periodic solution of equation (1.6). The boundary value problem for periodic solutions of (1.6) is given in (3.15). This setting represents a control problem where the role of time is taken by the spatial variable y and the control $\delta(t)$ is used to determine the initial conditions $U(t, -1) = 0$ and $U_y(t, -1) = \delta(t)$ which guarantee that $U(t, 1) = 0$. Periodic solutions are obtained numerically by expanding solutions with Fourier series in time and with Chebyshev series in space. Examples of periodic solutions computed numerically for $k = 1$ and $q = 11$ are portrayed in Figures 2 and 3.

It is important to remark that the Newton-Kantorovich argument based on the radii polynomial approach can be used to validate rigorously the full branch of steady states u_λ or the eigenfunctions associated to the linear elliptic operator $A(\lambda)$ (e.g. using the approach of [33, 14]). On the other hand, the validation of the branches of periodic solutions requires further investigations due to the lack of compactness of the inverse of the linear hyperbolic operator $L(\lambda)$. Indeed, while computer-assisted proofs have been used to validate the existence of periodic solutions in a nonlinear wave equations in [2] and in a nonlinear ill-posed Boussinesq equation (modelling shallow water waves) in [4], in our case the inverse of the linear hyperbolic operator $L(\lambda)$ lacks the necessary compactness to apply a similar approach.

The rest of the paper proceeds as follows. In Section 2, we first prove Theorem 1.1 by combining a Lyapunov-Schmidt reduction and the Crandall-Rabinowitz theorem. In Section 3, we compute numerically the steady states, the spectrum of the elliptic operator and the periodic solutions of equation (1.6). Finally, in Section 4 we present a computer-assisted proof of Theorem 1.2 that allows obtaining rigorous control over the steady state u_{λ^*} at the critical parameter value.

2. Existence of an infinite number of branches of periodic solutions.

2.1. Properties of the linear elliptic operator. In this section we analyze the properties of the steady states u_λ of equation (1.1) and properties of the eigenfunctions of the linear elliptic operator $A(\lambda)$ defined in (1.2).

PROPOSITION 2.1. *The steady state $u_\lambda(x)$ is even in x and increasing in the interval $x \in [0, \pi/2]$. Moreover, the steady state u_λ is decreasing as a function of λ .*

Proof. This follows from the maximum principle, for instance see the results in [13] and [16]. \square

Now we present properties of the eigenvalues and eigenfunctions of the linear elliptic operator $A(\lambda)$. By Sturm–Liouville theory the eigenvalues $\mu_k(\lambda)$ are simple for all k and $\lambda \in [0, \lambda_*]$. Furthermore, we can order the eigenvalues of $A(\lambda)$ by $0 \leq \mu_1(\lambda) < \mu_2(\lambda) < \dots$ such that the eigenfunction $v_k(x; \lambda)$ corresponding to eigenvalue $\mu_k(\lambda)$ has $k - 1$ simple zeros in $(-\pi/2, \pi/2)$.

Remark 2.2. The eigenvalues $\mu_k(\lambda)$ for $k = 1, 2, 3$ are computed numerically for the rescaled problem with $\lambda^* = (\pi/2)^{-2} \lambda_*$ in Section 3 (see Figure 6). Note that, while we do not perform this in the present paper, our numerical computations of the eigenvalues and eigenfunctions could be validated rigorously using a Newton-Kantorovich argument based on the radii polynomial approach similarly to the proof of Theorem 1.2 in Section 4.

PROPOSITION 2.3. *For $\lambda \in [0, \lambda_*]$ we have the following estimate for the eigenvalues $\mu_k(\lambda)$,*

$$(2.1) \quad 2\lambda \leq k^2 - \mu_k(\lambda) \leq 9\lambda .$$

Proof. Let $A, B : H_0^2 \subset L^2 \rightarrow L^2$ be the operators

$$A = -\partial_x^2 - B, \quad B \stackrel{\text{def}}{=} \frac{2\lambda}{(1 + u_\lambda(x))^3}.$$

Then we have that $m_\lambda |u|_{L^2}^2 \leq \langle Bu, u \rangle \leq M_\lambda |u|_{L^2}^2$, where

$$M_\lambda \stackrel{\text{def}}{=} 2\lambda \sup_{x \in [-\pi/2, \pi/2]} \frac{1}{(1 + u_\lambda(x))^3} \quad \text{and} \quad m_\lambda \stackrel{\text{def}}{=} 2\lambda \inf_{x \in [-\pi/2, \pi/2]} \frac{1}{(1 + u_\lambda(x))^3}.$$

The operator $-\partial_x^2 : H_0^2 \subset L^2 \rightarrow L^2$ has eigenvalues k^2 and eigenfunctions $\cos kx$ for $k \in 2\mathbb{N} + 1$ and $\sin kx$ for $k \in 2\mathbb{N}^+$. Since $-\partial_x^2 = A + B$, the Courant-Fischer-Weyl theorem implies that

$$\mu_k(\lambda) + m_\lambda \leq k^2 \leq \mu_k(\lambda) + M_\lambda.$$

By the properties of $u_\lambda(x)$, we have that $m_\lambda = 2\lambda$ and $M_\lambda = 2\lambda (1 + u_\lambda(0))^{-3}$. Since $u_{\lambda_*}(0)$ is estimated in Theorem 1.2 with $u_{\lambda_*}(0) > -0.38834671892$, then

$$(1 + u_{\lambda_*}(0))^{-3} < 4.5$$

and

$$M_\lambda = 2\lambda(1 + u_\lambda(0))^{-3} \leq 2\lambda(1 + u_{\lambda_*}(0))^{-3} \leq 9\lambda.$$

□

PROPOSITION 2.4. *The eigenvalues $\mu_k(\lambda)$ are monotonically decreasing for $\lambda \in [0, \lambda_*]$.*

Proof. Let $\lambda_1 < \lambda_2$. Since u_λ is decreasing in λ , then

$$\frac{\lambda_1}{(1 + u_{\lambda_1}(x))^3} < \frac{\lambda_2}{(1 + u_{\lambda_2}(x))^3}.$$

Thus the operator

$$C \stackrel{\text{def}}{=} \left(\frac{2\lambda_2}{(1 + u_{\lambda_2}(x))^3} - \frac{2\lambda_1}{(1 + u_{\lambda_1}(x))^3} \right) : H_0^2 \subset L^2 \rightarrow L^2,$$

is positive definite, that is there exists $c \geq 0$ such that $\langle Cu, u \rangle_{L^2} \geq c \langle u, u \rangle_{L^2}$ for all $u \in L^2$. Applying the Courant-Fischer-Weyl theorem to the operator $A(\lambda_1) = A(\lambda_2) + C : H_0^2 \rightarrow L^2$ we obtain that

□

$$\mu_k(\lambda_1) \geq \mu_k(\lambda_2) + c > \mu_k(\lambda_2).$$

PROPOSITION 2.5. *The eigenfunction $v_k(x; \lambda)$ corresponding to the eigenvalue $\mu_k(\lambda)$ of the linear elliptic operator $A(\lambda)$ satisfies*

$$v_k(x; \lambda) = (-1)^{k+1} v_k(-x; \lambda).$$

Proof. Since $u_\lambda(x)$ is even, then $v_k(-x; \lambda)$ is also an eigenfunction of $A(\lambda)$ for the eigenvalue $\mu_k(\lambda)$. Since the eigenvalues are simple, then the eigenfunctions are unique up to a scalar multiple, that is $v_k(x; \lambda) = \pm v_k(-x; \lambda)$. But $v_k(x; \lambda)$ has $k - 1$ simple zeros in $(-\pi/2, \pi/2)$, then $v_k(x; \lambda) = v_k(-x; \lambda)$ if k is odd and $v_k(x; \lambda) = -v_k(-x; \lambda)$ if k is even. □

In particular, we have that

$$v_k(x; 0) = \begin{cases} \cos(kx), & k \in 2\mathbb{N} + 1 \\ \sin(kx), & k \in 2\mathbb{N}^+. \end{cases}$$

2.2. Properties of the linear hyperbolic operator. In this section we analyze the properties of the spectrum for the linear hyperbolic operator

$$(2.2) \quad L(\lambda, \nu)u \stackrel{\text{def}}{=} (p/q)^2 u_{tt} + A(\lambda)u.$$

We define C_{sym}^2 as the subspace of 2π -periodic even functions $u(t, x)$ satisfying Dirichlet boundary conditions $u(\pm\pi/2) = 0$. Thus functions $u \in C_{sym}^2$ have the expansion

$$u(t, x) = \sum_{(j,k) \in \mathbb{N} \times \mathbb{N}^+} u_{j,k} \cos(jt) v_k(x; \lambda), \quad u_{j,k} \in \mathbb{R}.$$

By the estimates (2.1), the norm

$$|u|_s^2 \stackrel{\text{def}}{=} \sum_{(j,k) \in \mathbb{N} \times \mathbb{N}^+} |u_{j,k}|^2 (j^2 + k^2 + 1)^s,$$

is equivalent to the standard Sobolev norm. Thus the standard Sobolev space can be defined by

$$H^s(S^1 \times [-\pi/2, \pi/2]; \mathbb{R}) \stackrel{\text{def}}{=} \{u(t, x) \in L^2 : |u|_s < \infty\}.$$

DEFINITION 2.6. *Since the embedding $H^s \subset C^2$ holds for $s \geq 3$, then the subspace*

$$(2.3) \quad H_{sym}^s \stackrel{\text{def}}{=} \{u(t, x) \in H^s : u(t, \pm\pi/2) = 0, \quad u(t, x) = u(-t, x)\}$$

is well defined for $s \geq 3$.

The linear hyperbolic map

$$L(\lambda, p/q) = (p/q)^2 \partial_t^2 + A(\lambda) : D(L) \subset H_{sym}^s \rightarrow H_{sym}^s$$

is a closed operator, where the subspace $D(L) \subset H_{sym}^s$ is closed under the norm

$$|u|_L^2 = |Lu|_s^2 + |u|_s^2.$$

The linear map L has eigenvalues

$$\mu_{j,k}(\lambda, p/q) \stackrel{\text{def}}{=} -(pj/q)^2 + \mu_k(\lambda)$$

and eigenfunctions $\cos(jt)v_k(x)$ for $(j, k) \in \mathbb{N} \times \mathbb{N}^+$.

LEMMA 2.7. *Let $\lambda \in [\varepsilon, \lambda_*]$. If $pj/q = k$, then $|\mu_{j,k}(\lambda, p/q)| \geq 2\varepsilon$. If $pj/q \neq k$, then*

$$|\mu_{j,k}(\lambda, p/q)| \geq (pj/q + k)/q - 9\lambda.$$

Proof. If $pj/q = k$, then

$$|\mu_{j,k}(\lambda, p/q)| = k^2 - \mu_k(\lambda) \geq 2\lambda \geq 2\varepsilon.$$

If $-pj/q + k \neq 0$, then $|-pj/q + k| \geq 1/q$ and

$$\left| -(pj/q)^2 + k^2 \right| \geq |-pj/q + k| |pj/q + k| \geq (pj/q + k)/q.$$

Thus we have

$$|\mu_{j,k}(\lambda, p/q)| \geq \left| -(pj/q)^2 + k^2 \right| - |\mu_k(\lambda) - k^2| \geq (pj/q + k)/q - 9\lambda. \quad \square$$

We define $\mathcal{N}(\lambda, p/q)$ as the set of lattice points where $L(\lambda, p/q)$ has a zero eigenvalue,

$$(2.4) \quad \mathcal{N}(\lambda, p/q) \stackrel{\text{def}}{=} \{(j, k) : \mu_{j,k}(\lambda, p/q) = 0\}.$$

DEFINITION 2.8. *Notation $b \lesssim a$ means that there is a positive constant C such that $b \leq Ca$.*

PROPOSITION 2.9. *If $\lambda \in [\varepsilon, \lambda_* - \varepsilon]$, then $|\mu_{j,k}(\lambda, p/q)| \gtrsim \varepsilon$ for any $(j, k) \notin \mathcal{N}$ and \mathcal{N} is a bounded set with*

$$(2.5) \quad \mathcal{N} \subset \{(j, k) \in \mathbb{N}^+ \times \mathbb{N}^+ : pj + qk < 9q^2\lambda_* + 1\}.$$

Proof. Since the first eigenvalue $\mu_1(\lambda)$ is positive for $\lambda \in [0, \lambda_*)$ and $\mu_1(\lambda_*) = 0$, then $\mu_{0,k}(\lambda, p/q) = \mu_k(\lambda) \gtrsim \varepsilon$ for $\lambda < \lambda_* - \varepsilon$, which implies that $(j, k) \in \mathcal{N}$ only if $j > 0$. From Lemma 2.7, we have that $|\mu_{j,k}(\lambda, p/q)| \gtrsim \varepsilon$ for any $(j, k) \in \mathbb{N}^+ \times \mathbb{N}^+$ except when $(pj/q + k)/q - 9\lambda < 1/q^2$. This inequality is equivalent to

$$pj + qk < 9\lambda q^2 + 1 < 9\lambda_* q^2 + 1.$$

□

2.3. The Lyapunov-Schmidt reduction. The $2\pi q/p$ -periodic solutions of equation (1.1) of the form $U(t, x) = u_\lambda(x) + u(pt/q, x)$, where $u(t, x)$ is an even 2π -periodic perturbation satisfying Dirichlet boundary conditions $u(t, \pm\pi/2) = 0$, are solutions of the equation

$$(p/q)^2 u_{tt} - u_{xx} + \frac{\lambda}{(1 + u_\lambda + u)^2} - \frac{\lambda}{(1 + u_\lambda)^2} = 0.$$

Thus the equation for the perturbation u from the steady state u_λ reads

$$(2.6) \quad L(\lambda, p/q)u + g(u) = 0,$$

where L is defined in (2.2) and the quadratic nonlinear operator g is

$$(2.7) \quad \begin{aligned} g(u; \lambda) &\stackrel{\text{def}}{=} \frac{\lambda}{(1 + u_\lambda + u)^2} - \frac{\lambda}{(1 + u_\lambda)^2} + \frac{2\lambda}{(1 + u_\lambda)^3}u \\ &= \lambda \frac{3(1 + u_\lambda) + 2u}{(1 + u_\lambda)^3 (1 + u_\lambda + u)^2} u^2. \end{aligned}$$

In this section we make a Lyapunov-Schmidt reduction; namely, we solve the equation (2.6) in the range of the operator L (the range equation) and we obtain an equivalent equation to (2.6) defined in the kernel of L (the bifurcation equation).

We start by defining the projection in the kernel of $L(\lambda_0)$ as

$$Qu = \sum_{(j,k) \in \mathcal{N}} u_{j,k} \cos(jt) v_k(x; \lambda_0) : H_{sym}^s \rightarrow H_{sym}^s,$$

and the projection $P = I - Q : H_{sym}^s \rightarrow H_{sym}^s$ in the complement to the kernel of $L(\lambda_0)$. We have the following result.

COROLLARY 2.10. *Since $|\mu_{j,k}(\lambda)|^{-1} \lesssim \varepsilon^{-1}$ for $(j, k) \notin \mathcal{N}$, then*

$$(2.8) \quad \left| (PLP)^{-1} u \right|_s \lesssim \varepsilon^{-1} |Pu|_s.$$

Thus $(PLP)^{-1} : PH_{sym}^s \rightarrow PH_{sym}^s$ is a bounded operator. However, the operator $(PLP)^{-1}$ is not compact because the embedding $D(L) \subset H^s$ is not necessarily compact.

For $s \geq 3$, the space H_{sym}^s is a Banach algebra. Since $g(u) = \mathcal{O}(|u|^2)$ is analytic in a neighborhood of $u = 0$, then the nonlinear operator

$$g(u) = \mathcal{O}(|u|_s^2) : \mathcal{B}_r \subset H_{sym}^s \rightarrow H_{sym}^s$$

is continuous and well defined in

$$\mathcal{B}_r \stackrel{\text{def}}{=} \{u \in H_{sym}^s : |u|_s < r\}.$$

Setting

$$v = Qu, \quad w = Pu,$$

then $u = v + w$. Thus solutions to equation (2.6) are solutions of the kernel equation

$$(2.9) \quad QLQv + Qg(v + w) = 0,$$

and the range equation

$$(2.10) \quad PLPw + Pg(v + w) = 0.$$

The solutions to the range equation (2.10) in a neighborhood of $(v, \lambda) = (0, \lambda_0)$ are fixed points of the operator

$$Kw \stackrel{\text{def}}{=} (PLP)^{-1} g(v + w) : P\mathcal{B}_r \subset H_{sym}^s \rightarrow H_{sym}^s.$$

The estimate (2.8) and the fact that $g(v + w) = \mathcal{O}(|w|_s^2)$ imply that K is a contraction from the domain $P\mathcal{B}_r$ into itself when we choose $r \ll \varepsilon$. By the contracting mapping theorem there is a unique fixed point $w(v, \lambda) \in H_{sym}^s$ of K for (v, λ) in a neighborhood of $(0, \lambda_0)$. Thus there is a unique function $w(v, \lambda) \in H_{sym}^s$ that solves the range equation (2.10) in a neighborhood of $(0, \lambda_0)$. We conclude that the solutions to the equation (2.6) are given by the solutions of the bifurcation equation

$$(2.11) \quad QLQv + Qg(v + w(v, \lambda)) : Q\mathcal{B}_r \subset \ker L(\lambda_0) \rightarrow \ker L(\lambda_0).$$

2.4. The bifurcation equation. To solve the bifurcation equation (2.11) we need to look for values $\lambda_0 \in (0, \lambda_*)$ such that the linearization $L(\lambda_0)$ has a nontrivial kernel, that is $\mu_{j,k}(\lambda_0, p/q) = 0$ for some lattice point $(j, k) \in \mathcal{N}$.

PROPOSITION 2.11. *Define*

$$B_k \stackrel{\text{def}}{=} \left(\sqrt{\mu_k(\lambda_*)}, k \right) \subset \left(\sqrt{k^2 - 2\lambda_*}, k \right).$$

For each rational $p/q \in B_k$ there is a **unique** $\lambda_0 \in (0, \lambda_*)$ such that $\mu_{1,k}(\lambda_0, p/q) = 0$.

Proof. Since $\mu_k(\lambda; p/q)$ is decreasing and continuous for $\lambda \in (0, \lambda_*)$, any λ_0 such that $\mu_{1,k}(\lambda_0, p/q) = -(p/q)^2 + \mu_k(\lambda_0) = 0$ is unique. The result follows from the fact that the eigenvalue $\mu_k(\lambda)$ goes from the value $\mu_k(\lambda_*) \leq k^2 - 2\lambda_*$ to $\mu_k(0) = k^2$ for λ in the interval $(0, \lambda_*)$. \square

For each $\lambda_0 \in (0, \lambda_*)$ such that $\mu_{j,k}(\lambda_0, p/q) = 0$, the set $\mathcal{N}(\lambda_0, p/q)$ representing the kernel of $L(\lambda_0, p/q)$ may contain additional resonant points. If these resonances exist, they are contained in the bounded set given in (2.5).

DEFINITION 2.12. *We say that λ_0 is a non-resonant value if*

$$(2.12) \quad \mathcal{N}(\lambda_0, p/q) = \{(j, k) \in \mathbb{N} \times \mathbb{N}^+ : \mu_{j,k}(\lambda_0, p/q) = 0\} = \{(1, k)\}.$$

If λ_0 is non-resonant, then the kernel has dimension one, that is

$$\ker L(\lambda_0, p/q) = \{b \cos(t)v_k(x; \lambda_0) : b \in \mathbb{R}\}.$$

To prove the existence of a simple bifurcation, we need to choose non-resonant values λ_0 . The following lemma assures the existence of an infinite set of non-resonant values λ_0 .

PROPOSITION 2.13. *The set of non-resonant points $\lambda_0 \in (0, \lambda_*)$ such that*

$$\mu_{1,k}(\lambda_0, p/q) = 0$$

for some $p/q \in B_k$ is infinite.

Proof. There is a dense set of rationals $p/q \in B_k$ such that $\mu_{1,k}(\lambda_0, p/q) = 0$ for some $\lambda_0 \in (0, \lambda_*)$. Fix one of those points λ_0 . By (2.5), there is at most a finite number of resonant elements $(j_m, k_m) \in \mathbb{N}^+ \times \mathbb{N}^+$ such that $\mu_{j_m, k_m}(\lambda_0, p/q) = -(pj_m/q)^2 + \mu_{k_m}(\lambda_0) = 0$ for $m \in \{0, \dots, M\}$. Since the eigenvalues $\mu_k(\lambda_0)$ are simple, then

$$(pj_m/q)^2 = \mu_{k_m}(\lambda_0) \neq \mu_{k_0}(\lambda_0) = (pj_0/q)^2.$$

Therefore, the numbers j_m 's are different for different numbers $m \in \{0, \dots, M\}$. This implies that there is a unique lattice point denoted by (j_0, k_0) such that j_0 is maximal, that is $j_m < j_0$ for $m = 1, \dots, M$.

By choosing $p_0 = pj_0$ we have that $\mu_{1, k_0}(\lambda_0, p_0/q) = 0$ and $\mu_{j, k}(\lambda_0, p_0/q) \neq 0$ for all $j > 1$; otherwise j_0 would not be maximal. We conclude that any λ_0 such that $\ker L(\lambda_0, p_0/q)$ is not trivial is non-resonant for a rational number $p_0/q \in B_{k_0}$, that is $\mathcal{N}(\lambda_0, p_0/q) = \{(1, k_0)\}$. Moreover, the set of non-resonant values λ_0 for rational numbers $p_0/q \in B_{k_0}$ is infinite, otherwise there has to be at least one point $\lambda_0 \in (0, \lambda_*)$ with an infinite number of resonances, which is a contradiction to (2.5). \square

Remark 2.14. The choice of the maximal p_0 is equivalent to the choice of the minimal period $T = 2\pi q/p_0$. This argument is similar to the argument used in [17], [6] and [15].

PROPOSITION 2.15. *If λ_0 is a non-resonant value with $\mu_{1, k_0}(\lambda_0) = 0$, then equation (2.6) has a local bifurcation of zeros from $(u, \lambda_0) = (0, \lambda_0)$ such that*

$$(2.13) \quad u(x, t) = b(\cos t) v_{k_0}(x) + \mathcal{O}_{C_{sym}^2}(b^2), \quad \lambda = \lambda_0 + \mathcal{O}(b),$$

where $b \in [0, b_0]$ and $\mathcal{O}_{C_{sym}^2}(b^2)$ is a function in C_{sym}^2 of order b^2 .

Proof. The estimate

$$v = b(\cos t) v_{k_0}(x) + \mathcal{O}(b^2), \quad \lambda = \lambda_0 + \mathcal{O}(b),$$

is a consequence of the Crandall-Rabinowitz theorem. Thus the result follows from the Lyapunov-Schmidt reduction and the fact that $w(v; \lambda) = \mathcal{O}_{H_{sym}^s}(v^2)$ with $H_{sym}^s \subset C_{sym}^2$ for $s \geq 3$. To apply the Crandall-Rabinowitz theorem, we only need to verify that $\partial_\lambda L(\lambda_0)(\cos t) v_{k_0}(x)$ is not in the range of $L(\lambda_0)$,

$$\langle \partial_\lambda L(\lambda_0)(\cos t) v_{k_0}(x), (\cos t) v_{k_0}(x) \rangle_{L^2} \neq 0.$$

This condition is equivalent to

$$\int_{-\pi/2}^{\pi/2} \left(\partial_\lambda \frac{\lambda}{(1+u_\lambda)^3} \right)_{\lambda=\lambda_0} v_{k_0}^2(x) dx \neq 0,$$

which follows from the fact that u_λ is decreasing in λ . That is, we have $-\partial_\lambda u_\lambda \geq 0$

and

$$\left(\partial_\lambda \frac{\lambda}{(1+u_\lambda)^3} \right)_{\lambda=\lambda_0} = \frac{1}{(1+u_\lambda)^3} - \partial_\lambda u_\lambda \frac{\lambda}{(1+u_\lambda)^4} > 0. \quad \square$$

2.5. Proof of Theorem 1.1. The proof of Theorem 1.1 is a consequence of Proposition 2.15 and the fact that

$$U(t, x) = u_\lambda(x) + u(pt/q, x).$$

It only remains to obtain the symmetries of the local bifurcations. This is a consequence of the following proposition.

PROPOSITION 2.16. *The bifurcation (2.13) has the symmetries $u(t, x) = u(t, -x)$ for k odd and $u(t, x) = u(t + \pi, -x)$ for k even.*

Proof. Since $u_\lambda(x)$ is even, the equation

$$L(\lambda)u + g(u) = 0$$

is equivariant under the action of the group $(\kappa_1, \kappa_2) \in \mathbb{Z}_2 \times \mathbb{Z}_2$ in $u(t, x) \in H_{sym}^s$ given by

$$\kappa_1 u(t, x) = u(t, -x), \quad \kappa_2 u(t, x) = u(t + \pi, x).$$

Since

$$u(t, x) = \sum_{(j,k) \in \mathbb{N} \times \mathbb{N}^+} u_{j,k} \cos(jt) v_k(x),$$

where $v_k(-x) = (-1)^{k+1} v_k(x)$ and $\cos j(t + \pi) = (-1)^j \cos jt$, the actions of κ_1 and κ_2 in the components $u_{j,k} \in \mathbb{R}$ are given by $\kappa_1 u_{j,k} = (-1)^{k+1} u_{j,k}$ and $\kappa_2 u_{j,k} = (-1)^j u_{j,k}$. In particular, for $j = 1$ we have

$$\kappa_1 u_{1,k} = (-1)^{k+1} u_{1,k}, \quad \kappa_2 u_{1,k} = -u_{1,k}.$$

Then $u_{1,k} \in \ker L(\lambda_0)$ is fixed by the action of κ_1 if k is odd and by $\kappa_1 \kappa_2$ if k is even. The result follows from the fact that u is fixed by the action of κ_1 if it satisfies that

$$u(t, x) = \kappa_1 u(t, x) = u(t, -x),$$

and by $\kappa_1 \kappa_2$ if

$$u(t, x) = \kappa_1 \kappa_2 u(t, x) = u(t + \pi, -x). \quad \square$$

Remark 2.17. Since the equation $L(\lambda, p/q)u + g(u; \lambda) = 0$ has a gradient structure and the eigenvalues of $L(\lambda)$ cross zero in the same direction, because they are decreasing in λ , then one can use Conley index to prove that every value λ_0 where the kernel of $L(\lambda_0)$ is not trivial is a bifurcation point [25]. Therefore, for every $p/q \in B_k$ there is a (possibly resonant) bifurcation value λ_0 such that $\mu_{j,k}(\lambda_0; p/q) = 0$. Thus the set of bifurcation values λ_0 is dense when considering all possible values of p and q . However, the result using Conley index does not guarantee that the bifurcation is a continuum satisfying estimates (2.13).

3. Numerical computation of branches of periodic solutions. In this section we compute numerically the steady states, the spectrum of the elliptic operator and the periodic solutions of the scaled equation (1.6). This is done in Sections 3.1, 3.2 and 3.3, respectively. For each of these problems, we introduce an infinite dimensional zero-finding problem whose solutions correspond to the wanted objects of interest. Then, a standard predictor-corrector numerical continuation method is applied to finite dimensional projections of each problem. The reason of introducing first the infinite dimensional formulation of the problems is twofold. First, it matches the formulation of Section 4 involved in the rigorous computation of the saddle-node bifurcation. Second, it would allow to have the proper formulation for possibly doing rigorous computations of periodic solutions in the future.

3.1. Continuation of the branch of steady states. The steady states of the MEMS equation (1.6) satisfy the nonlinear boundary value problem

$$(3.1) \quad U_{yy} = \frac{\lambda}{(1+U)^2}, \quad U(-1) = U(1) = 0.$$

To compute the solutions of (3.1), we first transform the equation into a differential equation with polynomial nonlinearities. Letting $u_1 \stackrel{\text{def}}{=} U$, $u_2 \stackrel{\text{def}}{=} U_y = u'_1$ and $u_3 \stackrel{\text{def}}{=} \frac{1}{1+u_1}$ yields that

$$\begin{aligned} u'_1 &= u_2 \\ u'_2 &= U_{yy} = \frac{\lambda}{(1+U)^2} = \lambda u_3^2 \\ u'_3 &= -\frac{1}{(1+u_1)^2} u'_1 = -u_2 u_3^2. \end{aligned}$$

The boundary conditions $u_1(-1) = u_1(1) = 0$ are appended. To fix the right condition for u_3 , we impose that $u_3(-1) = \frac{1}{1+u_1(-1)} = 1$. The problem of computing a solution $U(y)$ of the nonlinear non-polynomial equation (3.1) is then transformed into the polynomial boundary value problem

$$(3.2) \quad \begin{pmatrix} u'_1 \\ u'_2 \\ u'_3 \end{pmatrix} = \begin{pmatrix} u_2 \\ \lambda u_3^2 \\ -u_2 u_3^2 \end{pmatrix}, \quad \begin{pmatrix} u_1(-1) \\ u_2(-1) \\ u_3(-1) \end{pmatrix} = \begin{pmatrix} 0 \\ \delta \\ 1 \end{pmatrix}, \quad u_1(1) = 0,$$

where δ is the unknown initial velocity $U'(-1)$ which we will solve for. We expand solutions with Chebyshev series

$$u_j(y) = (a_j)_0 + 2 \sum_{k \geq 0} (a_j)_k T_k(y), \quad j = 1, 2, 3,$$

where $T_n : [-1, 1] \rightarrow \mathbb{R}$ ($n \geq 0$) are the Chebyshev polynomials.

Denote by $u = (u_1, u_2, u_3)$ and $\Psi(u) \in \mathbb{R}^3$ the right-hand side of the polynomial differential equation given in (3.2). Denote $a_j = ((a_j)_n)_{n \geq 0}$ for $j = 1, 2, 3$, and $a = (a_1, a_2, a_3)$. For each $j = 1, 2, 3$, the Chebyshev expansion of $\Psi_j(u(\cdot)) : [-1, 1] \rightarrow \mathbb{R}$ is given by

$$\Psi_j(u(y)) = (c_j)_0 + 2 \sum_{n \geq 0} (c_j)_n T_n(y), \quad j = 1, 2, 3, \quad y \in [-1, 1],$$

where

$$(3.3) \quad \begin{pmatrix} c_1 \\ c_2 \\ c_3 \end{pmatrix} = \begin{pmatrix} c_1(a) \\ c_2(\lambda, a) \\ c_3(a) \end{pmatrix} \stackrel{\text{def}}{=} \begin{pmatrix} a_2 \\ \lambda a_3^2 \\ -a_2 a_3^2 \end{pmatrix},$$

where $a_3^2 = a_3 * a_3$ and $a_2 a_3^2 = a_2 * a_3 * a_3$ are standard discrete convolutions.

For $j = 1, 2, 3$, let

$$(3.4) \quad (f_j^{(\text{eq})}(\lambda, \delta, a))_n \stackrel{\text{def}}{=} \begin{cases} (a_j)_0 + 2 \sum_{\ell=1}^{\infty} (-1)^\ell (a_j)_\ell - \alpha_j, & n = 0, \\ 2n(a_j)_n + (c_j(a))_{n+1} - (c_j(a))_{n-1}, & n \geq 1, \end{cases}$$

where

$$(3.5) \quad \alpha_j \stackrel{\text{def}}{=} \begin{cases} 0, & j = 1 \\ \delta, & j = 2. \\ 1, & j = 3 \end{cases}$$

Setting $f_j^{(\text{eq})} = ((f_j^{(\text{eq})})_n)_{n \geq 0}$ and $\eta(a_1) = (a_1)_0 + 2 \sum_{\ell=1}^{\infty} (a_1)_\ell$ (this is the Chebyshev expansion of the extra condition $u_1(1) = 0$), the resulting map to solve in the space of Chebyshev coefficients is given by

$$(3.6) \quad f^{(\text{eq})}(\lambda, \delta, a) \stackrel{\text{def}}{=} \begin{pmatrix} \eta(a_1) \\ f_1^{(\text{eq})}(a) \\ f_2^{(\text{eq})}(\lambda, \delta, a) \\ f_3^{(\text{eq})}(a) \end{pmatrix}.$$

Define the operators (acting on Chebyshev sequences) by

$$(3.7) \quad T \stackrel{\text{def}}{=} \begin{pmatrix} 0 & 0 & 0 & 0 & 0 & \cdots \\ -1 & 0 & 1 & 0 & \cdots & \\ 0 & -1 & 0 & 1 & 0 & \cdots \\ & \ddots & \ddots & \ddots & \ddots & \ddots \\ & \cdots & 0 & -1 & 0 & 1 \\ & & \cdots & \ddots & \ddots & \ddots \end{pmatrix},$$

and

$$(3.8) \quad \Lambda \stackrel{\text{def}}{=} \begin{pmatrix} 0 & 0 & 0 & 0 & 0 & \cdots \\ 0 & 2 & 0 & 0 & \cdots & \\ 0 & 0 & 4 & 0 & 0 & \cdots \\ & \ddots & \ddots & \ddots & \ddots & \ddots \\ & \cdots & 0 & 0 & 2\ell & 0 \\ & & \cdots & \ddots & \ddots & \ddots \end{pmatrix}.$$

Using the above operators, we may write for the cases $n > 0$

$$(f_j^{(\text{eq})}(\lambda, \delta, a))_n = (\Lambda a_j + T c_j(a))_n = 2n(a_j)_n + (c_j(a))_{n+1} - (c_j(a))_{n-1}.$$

Hence, for $j = 1, 2, 3$,

$$(f_j^{(\text{eq})}(\lambda, \delta, a))_n \stackrel{\text{def}}{=} \begin{cases} (a_j)_0 + 2 \sum_{\ell=1}^{\infty} (-1)^\ell (a_j)_\ell - \alpha_j, & n = 0, \\ (\Lambda a_j + T c_j(a))_n, & n > 0. \end{cases}$$

By construction, computing solutions to the nonlinear BVP (3.1) (that is computing equilibria of the MEMS equation) boils down to computing simultaneously λ, δ, a such that $f^{(\text{eq})}(\lambda, \delta, a) = 0$ where $f^{(\text{eq})}$ is defined in (3.6). Letting $x \stackrel{\text{def}}{=} (\delta, a)$ we can compute branches of steady states by applying a parameter continuation method (that is a predictor corrector algorithm, see [23]) to a finite dimensional projection of the problem $f^{(\text{eq})}(\lambda, x) = 0$, where λ is a continuation parameter. Figure 2 contains the image of the stable branch of steady states (black branch) computed numerically using the presented method.

3.2. Continuation of eigenfunctions and eigenvalues. The eigenfunctions and eigenvalues of the linearized problem can be computed similarly. The eigenfunctions and eigenvalues are needed in order to find the initial predictor to compute numerically the branches of periodic solutions. The eigenvalue problem associated to (3.1) is given by

$$\begin{aligned} -U_{yy} + \frac{\lambda}{(1+U)^2} &= 0, & U(-1) = U(1) &= 0 \\ -V_{yy} - \frac{2\lambda}{(1+U)^3}V - \mu V &= 0, & V(-1) = V(1) &= 0, \end{aligned}$$

where (μ, V) is an eigenvalue-eigenvector couple associated to the linearization of the MEMS equation (1.6) about the steady state solution U .

Letting $u_1 = U$, $u_2 = U_y = u'_1$, $u_3 = \frac{1}{1+u_1}$, $u_4 = V$ and $u_5 = V_y$ yields the system

$$(3.9) \quad \begin{pmatrix} u'_1 \\ u'_2 \\ u'_3 \\ u'_4 \\ u'_5 \end{pmatrix} = \begin{pmatrix} u_2 \\ \lambda u_3^2 \\ -u_2 u_3^2 \\ u_5 \\ -2\lambda u_3^2 u_4 - \mu u_4 \end{pmatrix}, \quad \begin{pmatrix} u_1(-1) \\ u_2(-1) \\ u_3(-1) \\ u_4(-1) \\ u_5(-1) \end{pmatrix} = \begin{pmatrix} 0 \\ \delta_1 \\ 1 \\ 0 \\ \delta_2 \end{pmatrix}, \quad \begin{pmatrix} u_1(1) \\ u_4(1) \end{pmatrix} = \begin{pmatrix} 0 \\ 0 \end{pmatrix},$$

where δ_1 and δ_2 are to be (uniquely) determined. Denote $\delta = (\delta_1, \delta_2)$ and $u = (u_1, \dots, u_5)$. The unknown variables in the polynomial boundary value problem (3.9) are (δ, μ, u) . An extra phase condition (that is one which fixes the length of the eigenvector $u_4 = V$) will be imposed to isolate the solutions (and therefore allowing the use of Newton's method).

We solve the eigenvalue problem (3.9) using Chebyshev series expansion, similarly to the BVP (3.2). We expand solutions with Chebyshev series

$$u_j(y) = (a_j)_0 + 2 \sum_{k \geq 0} (a_j)_k T_k(y), \quad j = 1, \dots, 5.$$

Denote by $\Psi(u) \in \mathbb{R}^5$ the right-hand side of the polynomial differential equation given in (3.9). Denote $a_j = ((a_j)_n)_{n \geq 0}$ for $j = 1, \dots, 5$ and $a = (a_1, \dots, a_5)$. Assume that the Chebyshev expansion of $\Psi(u(y))$ is given by

$$\Psi(u(y))_j = (c_j)_0 + 2 \sum_{n \geq 0} (c_j)_n T_n(y), \quad j = 1, \dots, 5,$$

where $c = c(a) = (c_1, \dots, c_5)$ is given component-wise by

$$\begin{pmatrix} c_1 \\ c_2 \\ c_3 \\ c_4 \\ c_5 \end{pmatrix} \stackrel{\text{def}}{=} \begin{pmatrix} a_2 \\ \lambda a_3^2 \\ -a_2 a_3^2 \\ a_5 \\ -2\lambda a_3^2 a_4 - \mu a_4 \end{pmatrix},$$

where $a_3^2 = a_3 * a_3$, $a_2 a_3^2 = a_2 * a_3 * a_3$ and $a_3^2 a_4 = a_3 * a_3 * a_3 * a_4$ are discrete convolutions.

Denote $x \stackrel{\text{def}}{=} (\delta, \mu, a)$. For $j = 1, \dots, 5$, let

$$(3.10) \quad (g_j)_n(x) \stackrel{\text{def}}{=} \begin{cases} (a_j)_0 + 2 \sum_{\ell=1}^{\infty} (-1)^\ell (a_j)_\ell - \alpha_j, & n = 0, \\ 2n(a_j)_n + (c_j)_{n+1} - (c_j)_{n-1}, & n \geq 1, \end{cases}$$

where

$$(3.11) \quad \alpha_j \stackrel{\text{def}}{=} \begin{cases} 0, & j = 1 \\ \delta_1, & j = 2 \\ 1, & j = 3 \\ 0, & j = 4 \\ \delta_2, & j = 5 \end{cases}.$$

For $j = 1, \dots, 5$, we set $g_j \stackrel{\text{def}}{=} ((g_j)_n)_{n \geq 0}$. Let

$$\begin{aligned} \eta_1(a_1) &\stackrel{\text{def}}{=} (a_1)_0 + 2 \sum_{n=1}^{\infty} (a_1)_n \\ \eta_2(a_4) &\stackrel{\text{def}}{=} (a_4)_0 + 2 \sum_{n=1}^{\infty} (a_4)_n \\ \eta_3(a_4) &\stackrel{\text{def}}{=} l(a_4) - 1, \end{aligned}$$

where l is linear and acting as a phase condition for the eigenvector V (by fixing its length). Set $\eta \stackrel{\text{def}}{=} (\eta_1, \eta_2, \eta_3) \in \mathbb{R}^3$. The three extra conditions $\eta(a) = 0 \in \mathbb{R}^3$ are the extra conditions (in Chebyshev) enforcing that $u_1(1) = 0$, $u_4(1) = 0$ and that the eigenvector $u_4 = V$ is locally isolated. The resulting map to solve in the space of Chebyshev coefficients is given by

$$(3.12) \quad f^{(\text{lin})}(x, \lambda) \stackrel{\text{def}}{=} \begin{pmatrix} \eta(a) \\ g(x, \lambda) \end{pmatrix}.$$

We can then apply a standard predictor-corrector method to continue the eigenvalues $\mu_k(\lambda)$ for $k = 1, 2, 3$ and for $\lambda \in [0, \lambda^*)$. Having fixed $k \in \{1, 2, 3\}$, we begin the continuation at $\lambda = 0$ knowing theoretically that at the steady state $U = 0$, the eigenvalues are given by $\mu = \mu_k(0) = \left(\frac{\pi k}{2}\right)^2$, $k \geq 0$ with corresponding eigenvectors given by

$$(3.13) \quad V(y) = V_k(y) = \begin{cases} \sin\left(\frac{k\pi y}{2}\right), & k \text{ even,} \\ \cos\left(\frac{k\pi y}{2}\right), & k \text{ odd.} \end{cases}$$

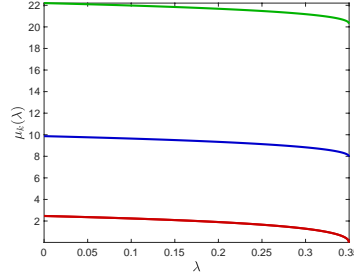


FIG. 6. Continuation of the eigenvalues $\mu_k(\lambda)$ for $k = 1, 2, 3$, $\lambda \in [0, \lambda^*)$.

Hence, recalling that $u_1 = U$, $u_2 = U_y = u_1'$, $u_3 = \frac{1}{1+u_1}$, $u_4 = V$ and $u_5 = V_y$, we get that at $\lambda = 0$, $u_1 \equiv 0$, $u_2 \equiv 0$, $u_3 \equiv 1$, $u_4(y) = V(y)$ as in (3.13) and

$$(3.14) \quad u_5(y) = V'(y) = V_k'(y) = \begin{cases} \frac{k\pi}{2} \cos\left(\frac{k\pi y}{2}\right), & k \text{ even,} \\ -\frac{k\pi}{2} \sin\left(\frac{k\pi y}{2}\right), & k \text{ odd.} \end{cases}$$

Using these explicit formulas for $u_j(y)$, $j = 1, \dots, 5$, we compute the corresponding Chebyshev series expansions to obtain the sequences $\bar{a}_1, \bar{a}_2, \dots, \bar{a}_5$. Note that the sequences $\bar{a}_1 = \bar{a}_2 = 0$ and $(\bar{a}_3)_n = \delta_{n,0}$, where $\delta_{i,j}$ is the Kronecker delta function. The computation of the Chebyshev coefficients \bar{a}_4 and \bar{a}_5 can be done analytically or using a numerical software. In our case, we use **Chebfun** to compute \bar{a}_4 and \bar{a}_5 . Moreover, we fix $\bar{\delta} = (\bar{\delta}_1, \bar{\delta}_2) = (u_2(-1), u_5(-1)) = (0, u_5(-1))$, where $u_5(-1)$ is determined exactly using (3.14). Letting $\bar{a} = (\bar{a}_1, \dots, \bar{a}_5)$, $\bar{\mu} = \left(\frac{\pi k}{2}\right)^2$, we have an approximate solution $\bar{x} = (\bar{\delta}, \bar{\mu}, \bar{a})$ which satisfies $f^{(\text{lin})}(\bar{x}, 0) \approx 0$. From that approximate solution at $\lambda = 0$ at a given $k \in \{1, 2, 3\}$, we perform a predictor-corrector continuation method on a finite dimensional projection of (3.12) to obtain a family of solutions of the form $\{x^{(k)}(\lambda) : \lambda \in [0, \lambda^*)\}$. Denote the second component of $x^{(k)}(\lambda)$ by $\mu_k(\lambda)$, we obtain a branch of eigenvalues parameterized over $\lambda \in [0, \lambda^*)$. See Figure 6 for a picture of the three branches $k = 1, 2, 3$.

3.3. Continuation of periodic solutions. The goal now is to compute periodic orbits $U = U(y, t)$ of (1.6). We fix two relatively prime integer (p, q) and fix a priori the frequency to be $\omega \stackrel{\text{def}}{=} \frac{\pi p}{2q}$. We aim at computing $\frac{2\pi}{\omega}$ -periodic orbits of (1.6). First, let us transform the problem into a polynomial one. Letting $U_1(y, t) \stackrel{\text{def}}{=} U(y, t)$, $U_2(y, t) \stackrel{\text{def}}{=} U_y(y, t) = \frac{d}{dy}U_1(y, t)$ and $U_3(y, t) \stackrel{\text{def}}{=} \frac{1}{1+U(y, t)} = \frac{1}{1+U_1(y, t)}$,

$$\begin{aligned} \frac{\partial}{\partial y}U_1 &= U_2 \\ \frac{\partial}{\partial y}U_2 &= U_{yy} = \frac{\lambda}{(1+U)^2} + U_{tt} = \lambda U_3^2 + \frac{\partial^2}{\partial t^2}U_1 \\ \frac{\partial}{\partial y}U_3 &= -\frac{1}{(1+U(y, t))^2}U_y(y, t) = -U_2U_3^2. \end{aligned}$$

The boundary conditions $U_1(-1, t) = U_1(1, t) = 0$ for all $t \in \mathbb{R}$ are appended. To fix the right condition for U_3 , we impose that $U_3(-1, t) = \frac{1}{1+U(-1, t)} = 1$ for all $t \in \mathbb{R}$.

After rescaling time from $[0, \frac{2\pi}{\omega}]$ to $[0, 2\pi]$, the problem of computing a solution $U(y, t)$ of the nonlinear hyperbolic equation (1.6) is therefore transformed into finding

a 2π -periodic orbit of the polynomial boundary value problem

$$(3.15) \quad \frac{\partial}{\partial y} \begin{pmatrix} U_1 \\ U_2 \\ U_3 \end{pmatrix} = \begin{pmatrix} U_2 \\ \lambda U_3^2 + \omega^2 \frac{\partial^2}{\partial t^2} U_1 \\ -U_2 U_3^2 \end{pmatrix}, \quad \begin{pmatrix} U_1(-1, t) \\ U_2(-1, t) \\ U_3(-1, t) \end{pmatrix} = \begin{pmatrix} 0 \\ \delta(t) \\ 1 \end{pmatrix}, \quad U_1(1, t) = 0,$$

for all $t \in \mathbb{R}$, where $\delta(t)$ is the a priori unknown initial velocity $U_y(-1, t)$ which we will solve for. We look for periodic orbits which are even in time. We expand U_j ($j = 1, 2, 3$) and δ in the form

$$(3.16) \quad U_j(y, t) = \sum_{\substack{n \geq 0 \\ k \geq 0}} (a_j)_{n,k} m_{n,k} T_n(y) \cos(kt) = \sum_{(n,k) \in \mathbb{Z}^2} (a_j)_{n,k} e^{i(n\theta + kt)}$$

$$(3.17) \quad \delta(t) = \sum_{k \geq 0} \delta_k m_{0,k} \cos(kt) = \sum_{k \in \mathbb{Z}} \delta_k e^{ikt}$$

where $\theta \stackrel{\text{def}}{=} \cos^{-1}(y)$,

$$m_{n,k} \stackrel{\text{def}}{=} \begin{cases} 1, & n = k = 0 \\ 2 & n = 0 \text{ and } k > 0 \\ 2 & n > 0 \text{ and } k = 0 \\ 4 & n \neq 0 \text{ and } k \neq 0, \end{cases}$$

and where $(a_j)_{n,k} = (a_1)_{|n|,|k|}$ and $\delta_{n,k} = \delta_{|n|,|k|}$ for $(n, k) \in \mathbb{Z}^2$. Denote $a_1 = ((a_1)_{n,k})_{n,k \geq 0}$, $a_2 = ((a_2)_{n,k})_{n,k \geq 0}$, $a_3 = ((a_3)_{n,k})_{n,k \geq 0}$ and $\delta = (\delta_k)_{k \geq 0}$. The unknowns are then given by

$$x \stackrel{\text{def}}{=} (\delta, a_1, a_2, a_3).$$

We integrate the BVP (3.15) in x to get the integral formulation

$$(3.18) \quad \begin{pmatrix} U_1(y, t) \\ U_2(y, t) \\ U_3(y, t) \end{pmatrix} = \begin{pmatrix} 0 \\ \delta(t) \\ 1 \end{pmatrix} + \int_{-1}^y \begin{pmatrix} U_2(\xi, t) \\ \lambda U_3^2(\xi, t) + \omega^2 \frac{\partial^2}{\partial t^2} U_1(\xi, t) \\ -U_2(\xi, t) U_3(\xi, t)^2 \end{pmatrix} d\xi, \quad y \in [-1, 1]$$

supplemented with the boundary condition $U_1(1, t) = 0$. We denote by $\Psi(U_1, U_2, U_3)$ the right-hand side of the polynomial problem (3.15). That is,

$$\Psi \begin{pmatrix} U_1(y, t) \\ U_2(y, t) \\ U_3(y, t) \end{pmatrix} \stackrel{\text{def}}{=} \begin{pmatrix} U_2 \\ \lambda U_3^2 + \omega^2 \frac{\partial^2}{\partial t^2} U_1 \\ -U_2 U_3^2 \end{pmatrix},$$

and for $j = 1, 2, 3$, we have the expansion

$$\Psi_j \begin{pmatrix} U_1(y, t) \\ U_2(y, t) \\ U_3(y, t) \end{pmatrix} = \sum_{\substack{n \geq 0 \\ k \geq 0}} (c_j)_{n,k} m_{n,k} T_n(y) \cos(kt) = \sum_{(n,k) \in \mathbb{Z}^2} (c_j)_{n,k} e^{i(n\theta + kt)},$$

where the terms c_1 , c_2 and c_3 involve discrete convolution terms. Explicitly,

$$\begin{aligned} (c_1)_{n,k} &= (a_2)_{n,k}, \\ (c_2)_{n,k} &= \lambda (a_3^2)_{n,k} - \omega^2 k^2 (a_1)_{n,k}, \\ (c_3)_{n,k} &= -(a_2 a_3^2)_{n,k}, \end{aligned}$$

where $a_3^2 = a_3 * a_3$ and $a_2 a_3^2 = a_2 * a_3 * a_3$ are standard two-dimensional discrete convolutions; for instance,

$$(a * b)_{n,k} = \sum_{\substack{n_1+n_2=n \\ k_1+k_2=k \\ k_i, n_i \in \mathbb{Z}}} a_{n_1, k_1} b_{n_2, k_2}.$$

For $j = 1, 2, 3$, let

$$(3.19) \quad (g_j)_{n,k}(x) \stackrel{\text{def}}{=} \begin{cases} (a_j)_{0,k} + 2 \sum_{\ell=1}^{\infty} (-1)^\ell (a_j)_{\ell,k} - (\alpha_j)_k, & n = 0, \\ 2n(a_j)_{n,k} + (c_j)_{n+1,k} - (c_j)_{n-1,k}, & n \geq 1, \end{cases}$$

where

$$(3.20) \quad (\alpha_j)_k \stackrel{\text{def}}{=} \begin{cases} 0, & j = 1 \\ \delta_k, & j = 2 \\ \hat{1}_k, & j = 3. \end{cases}$$

Setting $g_j = ((g_j)_{n,k})_{n,k \geq 0}$ and

$$\eta_k(a_1) \stackrel{\text{def}}{=} (a_1)_{0,k} + 2 \sum_{n \geq 1} (a_1)_{n,k} = 0,$$

which is the Fourier-Chebyshev expansion of the extra condition $U_1(1, t) = 0$ for all t . The resulting map to solve in Fourier-Chebyshev coefficients space is given by

$$(3.21) \quad f^{(\text{per})}(x, \lambda) \stackrel{\text{def}}{=} \begin{pmatrix} \eta(a_1) \\ g_1(x) \\ g_2(x, \lambda) \\ g_3(x) \end{pmatrix}.$$

Having identified a map whose zeros correspond to periodic orbits, we wish to compute (once more) branches of solutions (that is of periodic orbits) using a continuation method. The continuation requires first providing an initial point. Fix $k \geq 1$. For each rational λ_0 such that $\mu_k(\lambda_0) = \omega^2$ with $\omega = \frac{\pi p}{2q}$, there is a local continuum of 2π -periodic solution bifurcating from the steady solution $u_\lambda(y)$. The initial periodic orbit (that is the predictor) is given by

$$\hat{u}(t, y) \stackrel{\text{def}}{=} u_{\lambda_0}(y) + b \cos(t) v_k(y; \lambda_0),$$

for a small b . After having transformed this initial point as a sequence of Fourier-Chebyshev coefficients, we initiate the pseudo-arclength continuation (e.g. see [23]) on a finite dimensional projection of the map $f^{(\text{per})}$ defined in (3.21). Using that approach we performed several branch continuation, which are portrayed in Figures 2 and 3 for $k = 1$ and $q = 11$, and in Figures 4 and 5 for $k = 2$ and $q = 47$.

4. Rigorous computation of the saddle-node bifurcation. In this section, we prove Theorem 1.2. The proof of the theorem is computer-assisted and is based on the successful verification of the contraction mapping theorem of Newton-like operator acting on a ball of radius $r = 5.7 \times 10^{-12}$ centered at a numerical approximation of

a carefully chosen map defined on a Banach space X of fast decaying Chebyshev coefficients (the *saddle-node map* as defined in (4.6)).

To define the space X we require first to define the sequence space

$$(4.1) \quad \ell_\nu^1 \stackrel{\text{def}}{=} \left\{ \alpha = (\alpha_n)_{n \geq 0} : \|\alpha\|_\nu = |\alpha_0| + 2 \sum_{n \geq 1} |\alpha_n| \nu^n < \infty \right\},$$

for some fixed number $\nu \geq 1$. An important property of ℓ_ν^1 is that it is a Banach algebra under discrete convolutions, that is $\|a * b\|_\nu \leq \|a\|_\nu \|b\|_\nu$ for all $a, b \in \ell_\nu^1$. This is useful to perform the necessary estimates to analyze the nonlinear map. We also denote by

$$(4.2) \quad \tilde{\ell}_\nu^1 \stackrel{\text{def}}{=} \left\{ \alpha = (\alpha_n)_{n \geq 0} : \|\alpha\|_\nu = |\alpha_0| + 2 \sum_{n \geq 1} |\alpha_n| \frac{\nu^n}{n} < \infty \right\},$$

to the corresponding space of Chebyshev coefficients with slightly less decay (regularity) than ℓ_ν^1 . Note that for any $\nu \geq 1$ and for a fixed λ , one can show that the map $f^{(\text{eq})}$ defined in (3.6) satisfies

$$f^{(\text{eq})} : \mathbb{R} \times (\ell_\nu^1)^3 \rightarrow \mathbb{R} \times (\tilde{\ell}_\nu^1)^3.$$

In order to construct the saddle-node map F , we let

$$(4.3) \quad g(\lambda, a, \gamma, b) \stackrel{\text{def}}{=} D_{\delta, a} f^{(\text{eq})}(\lambda, \delta, a) \begin{pmatrix} \gamma \\ b \end{pmatrix} = \begin{pmatrix} \eta(b_1) \\ g_1(b) \\ g_2(\lambda, a, \gamma, b) \\ g_3(a, b) \end{pmatrix},$$

where

$$(g_j)_n \stackrel{\text{def}}{=} \begin{cases} (b_j)_0 + 2 \sum_{\ell=1}^{\infty} (-1)^\ell (b_j)_\ell - \tilde{\alpha}_j, & n = 0, \\ (\Lambda b_j + T d_j(a, b))_n, & n > 0, \end{cases}$$

for $j = 1, 2, 3$, with

$$(4.4) \quad \begin{pmatrix} d_1 \\ d_2 \\ d_3 \end{pmatrix} = \begin{pmatrix} d_1(b) \\ d_2(\lambda, a, b) \\ d_3(a, b) \end{pmatrix} \stackrel{\text{def}}{=} \begin{pmatrix} b_2 \\ 2\lambda a_3 b_3 \\ -a_3^2 b_2 - 2a_2 a_3 b_3 \end{pmatrix},$$

and

$$(4.5) \quad \tilde{\alpha}_j \stackrel{\text{def}}{=} \begin{cases} 0, & j = 1 \\ \gamma, & j = 2 \\ 0, & j = 3. \end{cases}$$

According to (3.6), computing the saddle-node bifurcation point (λ, δ, a) requires solving the augmented system

$$(4.6) \quad F(x) \stackrel{\text{def}}{=} \begin{pmatrix} \ell(\frac{\gamma}{b}) - 1 \\ f^{(\text{eq})}(\lambda, \delta, a) \\ D_{\delta, a} f^{(\text{eq})}(\lambda, \delta, a) \begin{pmatrix} \gamma \\ b \end{pmatrix} \end{pmatrix} = \begin{pmatrix} \ell(\frac{\gamma}{b}) - 1 \\ f^{(\text{eq})}(\lambda, \delta, a) \\ g(\lambda, a, \gamma, b) \end{pmatrix} = 0,$$

where

$$(4.7) \quad x \stackrel{\text{def}}{=} (\lambda, \delta, a, \gamma, b) \in X \stackrel{\text{def}}{=} \mathbb{R} \times \mathbb{R} \times (\ell_\nu^1)^3 \times \mathbb{R} \times (\ell_\nu^1)^3,$$

$b = (b_1, b_2, b_3) \in (\ell_\nu^1)^3$, and $\ell : \mathbb{R} \times (\ell_\nu^1)^3 \rightarrow \mathbb{R}$ is a linear functional acting on the eigenvector $\binom{\gamma}{b}$. We call the map F in (4.6) the *saddle-node map*. By construction, a non-degenerate zero \tilde{x} of F yields the existence of a saddle-node point, that is a point such that $f^{(\text{eq})}(\lambda, \delta, a) = 0$ and such that $D_{\delta, a} f^{(\text{eq})}(\lambda, \delta, a)$ has a one-dimensional kernel.

We endow the space X with the product norm

$$(4.8) \quad \|x\|_X \stackrel{\text{def}}{=} \max \{|\lambda|, |\delta|, \|a_1\|_\nu, \|a_2\|_\nu, \|a_3\|_\nu, |\gamma|, \|b_1\|_\nu, \|b_2\|_\nu, \|b_3\|_\nu\}.$$

Moreover, recalling (4.2), we define

$$(4.9) \quad Y \stackrel{\text{def}}{=} \mathbb{R} \times \mathbb{R} \times (\tilde{\ell}_\nu^1)^3 \times \mathbb{R} \times (\tilde{\ell}_\nu^1)^3,$$

and one can easily verify that $F : X \rightarrow Y$ is well defined.

In this section, we present a computer-assisted approach to solving the saddle-node map (4.6) using the tools of rigorous numerics in order to obtain a rigorous control over the value of λ^* . This approach will give a proof of Theorem 1.2. The idea of the computer-assisted proof is to demonstrate that a certain Newton-like operator is a contraction on a closed ball centered at a numerical approximation \bar{x} . To compute \bar{x} , we consider a finite dimensional projection of the saddle-node map $F : X \rightarrow Y$.

Given a number $m \in \mathbb{N}$, and given a vector $a = (a_n)_{n \geq 0} \in \ell_\nu^1$, consider the projection

$$\begin{aligned} \pi^m : \ell_\nu^1 &\rightarrow \mathbb{R}^m \\ a &\mapsto \pi^m a \stackrel{\text{def}}{=} (a_n)_{n=0}^{m-1} \in \mathbb{R}^m. \end{aligned}$$

Given $N \in \mathbb{N}$, we generalize that projection to get $\pi_N^m : (\ell_\nu^1)^N \rightarrow \mathbb{R}^{Nm}$ defined by

$$\pi_N^m(a^{(1)}, \dots, a^{(N)}) \stackrel{\text{def}}{=} (\pi^m a^{(1)}, \dots, \pi^m a^{(N)}) \in \mathbb{R}^{Nm},$$

and $\Pi^{(m)} : \mathcal{X} \rightarrow \mathbb{R}^{6m+3}$ defined by

$$\Pi^{(m)} x = \Pi^{(m)}(\lambda, \delta, a, \gamma, b) \stackrel{\text{def}}{=} (\lambda, \delta, \pi_3^m a, \gamma, \pi_3^m b) \in \mathbb{R}^{6m+3}.$$

Often, given $x \in X$, we denote

$$x^{(m)} \stackrel{\text{def}}{=} \Pi^{(m)} x \in \mathbb{R}^{6m+3}.$$

Moreover, we define the natural inclusion $\iota^m : \mathbb{R}^m \hookrightarrow \ell_\nu^1$ as follows. For $a = (a_n)_{n=0}^{m-1} \in \mathbb{R}^m$, we define $\iota^m a \in \ell_\nu^1$ component-wise by

$$(\iota^m a)_\ell = \begin{cases} a_n, & n = 0, \dots, m-1 \\ 0, & n \geq m. \end{cases}$$

Similarly, let $\iota_N^m : \mathbb{R}^{Nm} \hookrightarrow (\ell_\nu^1)^N$ be the natural inclusion defined as follows. Given $a = (a^{(1)}, \dots, a^{(N)}) \in (\mathbb{R}^m)^N \cong \mathbb{R}^{Nm}$, we define

$$\iota_N^m a \stackrel{\text{def}}{=} (\iota^m a^{(1)}, \dots, \iota^m a^{(N)}) \in (\ell_\nu^1)^N.$$

We define the natural inclusion $\iota^{(m)} : \mathbb{R}^{6m+3} \hookrightarrow X$, for $x \in \mathbb{R}^{6m+3}$, by

$$\iota^{(m)}x = \iota^{(m)}(\lambda, \delta, a, \gamma, b) \stackrel{\text{def}}{=} (\lambda, \delta, \iota_3^m a, \gamma, \iota_3^m b) \in X.$$

Let the *finite dimensional projection* $F^{(m)} : \mathbb{R}^{6m+3} \rightarrow \mathbb{R}^{6m+3}$ of the saddle-node map (4.6) be defined, for $x \in \mathbb{R}^{6m+3}$, as

$$(4.10) \quad F^{(m)}(x) = \Pi^{(m)}F(\iota^{(m)}x).$$

Assume that a numerical approximation $\bar{x} \in \mathbb{R}^{6m+3}$ of (4.10) has been obtained using Newton's method, that is $F^{(m)}(\bar{x}) \approx 0$. We slightly abuse the notation and denote $\bar{x} \in \mathbb{R}^{6m+3}$ and $\iota^{(m)}\bar{x} \in X$ both using \bar{x} .

The following result is a Newton-Kantorovich theorem with a smoothing approximate inverse. It provides an a-posteriori validation method for proving rigorously the existence of a point \tilde{x} such that $F(\tilde{x}) = 0$ and $\|\tilde{x} - \bar{x}\|_X \leq r$ for a small radius r . Recalling the norm on X given in (4.7), denote by

$$B_r(y) \stackrel{\text{def}}{=} \{x \in X : \|x - y\|_X \leq r\} \subset X$$

the ball of radius r centered at $y \in X$.

THEOREM 4.1 (Radii Polynomial Approach). *For $\bar{x} \in X$ and $r > 0$ assume that $F : X \rightarrow Y$ is Fréchet differentiable on the ball $B_r(\bar{x})$. Consider bounded linear operators $A^\dagger \in B(X, Y)$ and $A \in B(Y, X)$, where A^\dagger is an approximation of $DF(\bar{x})$ and A is an approximate inverse of $DF(\bar{x})$. Observe that*

$$(4.11) \quad AF : X \rightarrow X.$$

Assume that A is injective. Let $Y_0, Z_0, Z_1, Z_2 \geq 0$ be bounds satisfying

$$(4.12) \quad \|AF(\bar{x})\|_X \leq Y_0,$$

$$(4.13) \quad \|I - AA^\dagger\|_{B(X)} \leq Z_0,$$

$$(4.14) \quad \|A[DF(\bar{x}) - A^\dagger]\|_{B(X)} \leq Z_1,$$

$$(4.15) \quad \|A[DF(\bar{x} + z) - DF(\bar{x})]\|_{B(X)} \leq Z_2r, \quad \forall z \in B_r(0).$$

Define the radii polynomial

$$(4.16) \quad p(r) \stackrel{\text{def}}{=} Z_2r^2 + (Z_1 + Z_0 - 1)r + Y_0.$$

If there exists $0 < r_0 \leq r$ such that

$$(4.17) \quad p(r_0) < 0,$$

then there exists a unique $\tilde{x} \in B_{r_0}(\bar{x})$ such that $F(\tilde{x}) = 0$.

The proof of the theorem, which is a generalization of the usual Newton-Kantorovich theorem can be found (for example in [26]).

4.1. The operators A^\dagger and A . To apply the radii polynomial approach of Theorem 4.1, we need to define the operator A^\dagger (an approximation of the derivative $DF(\bar{x})$) and the operator A (an approximation of the inverse of $DF(\bar{x})$). Consider the finite dimensional projection $F^{(m)} : \mathbb{R}^{6m+3} \rightarrow \mathbb{R}^{6m+3}$ given in (4.10), and assume that we computed $\bar{x} \in \mathbb{R}^{6m+3}$ such that $F^{(m)}(\bar{x}) \approx 0$.

We denote by $DF^{(m)}(\bar{x}) \in M_{6m+3}(\mathbb{R})$ to the Jacobian matrix of $F^{(m)}$ at \bar{x} . For the sake of simplicity, given any $N \in \mathbb{N}$, we denote the differentiation operator D acting on $u \in (\ell_v^1)^N$ as

$$(4.18) \quad (Du)_n \stackrel{\text{def}}{=} 2nu_n = \begin{pmatrix} 2n(u_1)_n \\ 2n(u_2)_n \\ \vdots \\ 2n(u_N)_n \end{pmatrix}.$$

And given $x \in X$, we define

$$(4.19) \quad A^\dagger x = \iota^{(m)} \Pi^{(m)} A^\dagger x + (I - \iota^{(m)} \Pi^{(m)}) A^\dagger x,$$

where $\Pi^{(m)} A^\dagger x = DF^{(m)}(\bar{x})x^{(m)}$ and

$$(I - \iota^{(m)} \Pi^{(m)}) A^\dagger x = \begin{pmatrix} 0 \\ 0 \\ (I - \iota_3^m \pi_3^m) Da \\ 0 \\ (I - \iota_3^m \pi_3^m) Db \end{pmatrix}.$$

Recalling the definition of the Banach space Y in (4.9), we can verify that the operator $A^\dagger : X \rightarrow Y$ is a bounded linear operator. For m large enough, it acts as an approximation of the Fréchet derivative $D_x F(\bar{x})$. Its action on the finite dimensional projection is the Jacobian matrix (the derivative) of $F^{(m)}$ at \bar{x} , while its action on the tail only keeps the unbounded terms of the differentiation D defined in (4.18).

Now we consider a matrix $A^{(m)} \in M_{6m+3}(\mathbb{R})$ such that $A^{(m)} \approx DF^{(m)}(\bar{x})^{-1}$. In other words, this means that $\|I - A^{(m)} DF^{(m)}(\bar{x})\| \ll 1$. The computation of $A^{(m)}$ is done using a numerical software (MATLAB in our case). We decompose the matrix $A^{(m)}$ block-wise as

$$A^{(m)} = \begin{pmatrix} A_{\lambda,\lambda}^{(m)} & A_{\lambda,\delta}^{(m)} & A_{\lambda,a}^{(m)} & A_{\lambda,\gamma}^{(m)} & A_{\lambda,b}^{(m)} \\ A_{\delta,\lambda}^{(m)} & A_{\delta,\delta}^{(m)} & A_{\delta,a}^{(m)} & A_{\delta,\gamma}^{(m)} & A_{\delta,b}^{(m)} \\ A_{a,\lambda}^{(m)} & A_{a,\delta}^{(m)} & A_{a,a}^{(m)} & A_{a,\gamma}^{(m)} & A_{a,b}^{(m)} \\ A_{\gamma,\lambda}^{(m)} & A_{\gamma,\delta}^{(m)} & A_{\gamma,a}^{(m)} & A_{\gamma,\gamma}^{(m)} & A_{\gamma,b}^{(m)} \\ A_{b,\lambda}^{(m)} & A_{b,\delta}^{(m)} & A_{b,a}^{(m)} & A_{b,\gamma}^{(m)} & A_{b,b}^{(m)} \end{pmatrix},$$

so that it acts is defined on $x^{(m)} = (\lambda, \delta, a^{(m)}, \gamma, b^{(m)}) \in \mathbb{R}^{6m+3}$. Thus we define A as

$$(4.20) \quad A = \begin{pmatrix} A_{\lambda,\lambda} & A_{\lambda,\delta} & A_{\lambda,a} & A_{\lambda,\gamma} & A_{\lambda,b} \\ A_{\delta,\lambda} & A_{\delta,\delta} & A_{\delta,a} & A_{\delta,\gamma} & A_{\delta,b} \\ A_{a,\lambda} & A_{a,\delta} & A_{a,a} & A_{a,\gamma} & A_{a,b} \\ A_{\gamma,\lambda} & A_{\gamma,\delta} & A_{\gamma,a} & A_{\gamma,\gamma} & A_{\gamma,b} \\ A_{b,\lambda} & A_{b,\delta} & A_{b,a} & A_{b,\gamma} & A_{b,b} \end{pmatrix},$$

where the action of each block of A is finite (that is they act on $x^{(m)} = \Pi^{(m)}x$ only) except for the two diagonal blocks $A_{a,a}$ and $A_{b,b}$ which have infinite tails. More

explicitly, for each $j = 1, 2, 3$,

$$\begin{aligned} ((A_{a,a}a)_j)_n &= \begin{cases} ((A_{a,a}^{(m)}\pi_3^m a)_j)_n & \text{for } n = 0, \dots, m-1, \\ \frac{1}{2^n}(a_j)_n & \text{for } n \geq m, \end{cases} \\ ((A_{b,b}b)_j)_n &= \begin{cases} ((A_{b,b}^{(m)}\pi_3^m b)_j)_n & \text{for } n = 0, \dots, m-1, \\ \frac{1}{2^n}(b_j)_n & \text{for } n \geq m. \end{cases} \end{aligned}$$

Having defined the operators A and A^\dagger , we are ready to define the bounds Y_0 , Z_0 , Z_1 and Z_2 (satisfying (4.12), (4.13), (4.14) and (4.15), respectively), required to build the radii polynomial defined in (4.16).

4.2. Y_0 bound. Denote by $\bar{x} = (\bar{\lambda}, \bar{\delta}, \bar{a}, \bar{\gamma}, \bar{b}) \in X$ the numerical approximation with $\bar{a} = (\bar{a}_1, \bar{a}_2, \bar{a}_3) \in (\ell_\nu^1)^3$ and $\bar{b} = (\bar{b}_1, \bar{b}_2, \bar{b}_3) \in (\ell_\nu^1)^3$. Recalling the definition of f in (3.6) (which involves the convolutions in (3.3)) and the definition of g in (4.3) (which involves the convolutions in (4.4)), one has that

$$\begin{aligned} (I - \iota^{m+1}\pi^{m+1})f_1(\bar{a}) &= 0 \in \ell_\nu^1, \\ (I - \iota^{2m-1}\pi^{2m-1})f_2(\bar{\delta}, \bar{a}, \bar{\lambda}) &= 0 \in \ell_\nu^1, \\ (I - \iota^{3m-2}\pi^{3m-2})f_3(\bar{a}) &= 0 \in \ell_\nu^1, \\ (I - \iota^{m+1}\pi^{m+1})g_1(\bar{b}) &= 0 \in \ell_\nu^1, \\ (I - \iota^{2m-1}\pi^{2m-1})g_2(\bar{\lambda}, \bar{a}, \bar{\gamma}, \bar{b}) &= 0 \in \ell_\nu^1, \\ (I - \iota^{3m-2}\pi^{3m-2})g_3(\bar{a}, \bar{b}) &= 0 \in \ell_\nu^1. \end{aligned}$$

This result follows from the fact that the product of p trigonometric functions of degree $m-1$ is a trigonometric function of degree $p(m-1)$, and the n entry of the Chebyshev map f_j (resp. g_j) has entries of the form $(c_j)_{n+1} - (c_j)_{n-1}$ (resp. $(d_j)_{n+1} - (d_j)_{n-1}$). Using that information, one concludes that only finitely many entries of $F(\bar{x})$ are non-zeros, and therefore the computation of the bound Y_0 satisfying

$$(4.21) \quad \|AF(\bar{x})\|_X \leq Y_0$$

is a finite computation that can be performed using interval arithmetic (INTLAB in our case, see [30]).

4.3. Basic functional analytic background. In this section we present some elementary functional analytic background used to computing the bounds Z_0 and Z_1 .

For an infinite sequence of real numbers $a = \{a_n\}_{n \geq 0}$, and $\nu > 1$, we defined

$$\|a\|_\nu \stackrel{\text{def}}{=} |a_0| + 2 \sum_{n \geq 1} |a_n| \nu^n = \sum_{n \geq 0} |a_n| \omega_n,$$

where

$$\omega_n \stackrel{\text{def}}{=} \begin{cases} 1, & n = 0, \\ 2\nu^n, & n \geq 1. \end{cases}$$

The dual norm of $\|\cdot\|_\nu$ is

$$\|a\|_{\infty, \nu^{-1}} \stackrel{\text{def}}{=} \sup_{n \geq 0} \frac{|a_n|}{\omega_n},$$

and the set

$$\ell_{\infty, \nu-1} \stackrel{\text{def}}{=} \{ \{a_n\}_{n \geq 0} : \|a\|_{\infty, \nu-1} < \infty \},$$

is a Banach space.

We have the following known results (e.g. see [20]),

LEMMA 4.2. *If $a \in \ell_{\nu}^1$ and $c \in \ell_{\nu-1}^{\infty}$, then*

$$\left| \sum_{n \geq 0} c_n a_n \right| \leq \|c\|_{\infty, \nu-1} \|a\|_{\nu}.$$

LEMMA 4.3. *Given $\nu \geq 1$, $k \in \mathbb{Z}$ and $a \in \ell_{\nu}^1$, the function $l_a^k: \ell_{\nu}^1 \rightarrow \mathbb{C}$ defined by*

$$l_a^k(h) \stackrel{\text{def}}{=} (a * h)_k = \sum_{k_1+k_2=k} a_{k_1} h_{k_2}, \quad h \in \ell_{\nu}^1,$$

is a bounded linear functional, and

$$(4.22) \quad \|l_a^k\| = \sup_{\|h\|_{\nu} \leq 1} |l_a^k(h)| \leq \sup_{j \in \mathbb{Z}} \frac{|a_{k-j}|}{\nu^{|j|}} < \infty.$$

Fix a truncation mode to be m . Given $h \in \ell_{\nu}^1$, set

$$\begin{aligned} h^{(m)} &\stackrel{\text{def}}{=} (h_0, h_1, \dots, h_{m-1}, 0, 0, \dots) \in \ell_{\nu}^1, \\ h^{(\infty)} &\stackrel{\text{def}}{=} h - h^{(m)} \in \ell_{\nu}^1. \end{aligned}$$

COROLLARY 4.4. *Let $N \in \mathbb{N}$ and let $\bar{\alpha} = (\bar{\alpha}_0, \bar{\alpha}_1, \dots, \bar{\alpha}_N, 0, 0, \dots) \in \ell_{\nu}^1$. Suppose that $0 \leq k < m$ and define $\hat{l}_{\bar{\alpha}}^k \in \ell_{\nu-1}^{\infty}$ by*

$$\hat{l}_{\bar{\alpha}}^k(h) \stackrel{\text{def}}{=} (\bar{\alpha} * h^{(\infty)})_k = \sum_{k_1+k_2=k} \bar{\alpha}_{k_1} h_{k_2}^{(\infty)}.$$

Then, for all $h \in \ell_{\nu}^1$ such that $\|h\|_{\nu} \leq 1$,

$$(4.23) \quad \left| \hat{l}_{\bar{\alpha}}^k(h) \right| \leq \Psi_k(\bar{\alpha}) \stackrel{\text{def}}{=} \max_{j=m, \dots, k+m-1} \left(\frac{|\bar{\alpha}_{k-j} + \bar{\alpha}_{k+j}|}{2\nu^j} \right).$$

Proof. Notice that

$$\hat{l}_{\bar{\alpha}}^k(h) = (\bar{\alpha} * h^{(\infty)})_k = \sum_{|k_2| \geq m} \bar{\alpha}_{k-k_2} h_{k_2} = \sum_{k_2 \geq m} (\bar{\alpha}_{k-k_2} + \bar{\alpha}_{k+k_2}) h_{k_2} = \sum_{j \geq 0} c_j h_j,$$

where

$$c_j \stackrel{\text{def}}{=} \begin{cases} \bar{\alpha}_{k-j} + \bar{\alpha}_{k+j}, & \text{if } j \geq m, \\ 0, & \text{if } 0 \leq j \leq m. \end{cases}$$

Since $\bar{\alpha}_k = 0$ for all $|k| > N$, we use (4.22) to obtain

$$|l_{\bar{\alpha}}^k(h)| \leq \sup_{j \geq 0} \frac{|c_j|}{\omega_j} \leq \max_{j=m, \dots, k+m-1} \frac{|\bar{\alpha}_{k-j} + \bar{\alpha}_{k+j}|}{2\nu^j} = \Psi_k(\bar{\alpha}). \quad \square$$

4.4. Z_0 bound. We aim at computing a bound Z_0 satisfying (4.13). Let $B \stackrel{\text{def}}{=} I - AA^\dagger$, which we denote block-wise by

$$B = \begin{pmatrix} B_{\lambda,\lambda} & B_{\lambda,\delta} & B_{\lambda,a} & B_{\lambda,\gamma} & B_{\lambda,b} \\ B_{\delta,\lambda} & B_{\delta,\delta} & B_{\delta,a} & B_{\delta,\gamma} & B_{\delta,b} \\ B_{a,\lambda} & B_{a,\delta} & B_{a,a} & B_{a,\gamma} & B_{a,b} \\ B_{\gamma,\lambda} & B_{\gamma,\delta} & B_{\gamma,a} & B_{\gamma,\gamma} & B_{\gamma,b} \\ B_{b,\lambda} & B_{b,\delta} & B_{b,a} & B_{b,\gamma} & B_{b,b} \end{pmatrix}.$$

Note that the tails of B vanish by the definition of the diagonal tails of A and A^\dagger . We can compute the bound

$$Z_0^{(\alpha)} \stackrel{\text{def}}{=} \begin{cases} \sum_{\tilde{\alpha} \in \{\lambda, \delta, \gamma\}} |B_{\alpha, \tilde{\alpha}}| + \sum_{\substack{\tilde{\alpha} \in \{a_1, a_2, a_3, \\ b_1, b_2, b_3\}}} \|B_{\alpha, \tilde{\alpha}}\|_{\infty, \nu-1} & \alpha \in \{\lambda, \delta, \gamma\}, \\ \sum_{\substack{\tilde{\alpha} \in \{\lambda_1, \lambda_2, \lambda_3, \\ \alpha_1, \dots, \alpha_{n-1}\}}} \|B_{\alpha, \tilde{\alpha}}\|_{\nu} + \sum_{\substack{\tilde{\alpha} \in \{a_1, a_2, a_3, \\ b_1, b_2, b_3\}}} \max_{s=0, \dots, m-1} \frac{1}{\omega_s} \sum_{\ell=0, \dots, m-1} |(B_{\alpha, \tilde{\alpha}})_{\ell, s}| \omega_\ell & \alpha \in \{a_1, a_2, a_3, \\ b_1, b_2, b_3\}. \end{cases}$$

letting

$$(4.24) \quad Z_0 \stackrel{\text{def}}{=} \max \left\{ Z_0^{(\alpha)} : \alpha \in \{\lambda, \delta, \gamma, a_1, a_2, a_3, b_1, b_2, b_3\} \right\},$$

by construction we get that

$$\|I - AA^\dagger\|_{B(X)} \leq Z_0.$$

4.5. Z_1 bound. For any $h \in B_1(0) \subset X$, let

$$z \stackrel{\text{def}}{=} [DF(\bar{x}) - A^\dagger]h \in Y.$$

Denote

$$(4.25) \quad \begin{aligned} h &= (h_\lambda, h_\delta, h_{a_1}, h_{a_2}, h_{a_3}, h_\gamma, h_{b_1}, h_{b_2}, h_{b_3}), \\ z &= (z_\lambda, z_\delta, z_{a_1}, z_{a_2}, z_{a_3}, z_\gamma, z_{b_1}, z_{b_2}, z_{b_3}). \end{aligned}$$

Note that

$$z_\lambda = 0, \quad z_\delta = 2 \sum_{\ell \geq m} (h_{a_1})_\ell, \quad z_\gamma = 2 \sum_{\ell \geq m} (h_{b_1})_\ell.$$

Moreover, for $j = 1, 2, 3$, we have

$$(z_{a_j})_n = \begin{cases} 2 \sum_{\ell \geq m} (-1)^\ell (h_{a_j})_\ell, & n = 0 \\ (T\psi_{a_j})_n, & 0 < n \end{cases},$$

where

$$(\psi_{a_j})_n \stackrel{\text{def}}{=} \begin{cases} \begin{cases} 0, & 0 < n < m \\ (h_{a_2})_n, & n \geq m \end{cases}, & j = 1 \\ \begin{cases} (2\bar{\lambda}\bar{a}_3 h_{a_3}^{(\infty)})_n, & 0 < n < m \\ (2\bar{\lambda}\bar{a}_3 h_{a_3} + \bar{a}_3^2 h_\lambda)_n, & n \geq m \end{cases}, & j = 2 \\ \begin{cases} -(2\bar{a}_2\bar{a}_3 h_{a_3}^{(\infty)} + \bar{a}_3^2 h_{a_2}^{(\infty)})_n, & 0 < n < m \\ -(2\bar{a}_2\bar{a}_3 h_{a_3} + \bar{a}_3^2 h_{a_2})_n, & n \geq m \end{cases}, & j = 3. \end{cases}$$

Similarly, for $j = 1, 2, 3$, we have

$$(z_{b_j})_n = \begin{cases} 2 \sum_{\ell \geq m} (-1)^\ell (h_{b_j})_\ell, & n = 0 \\ (T\psi_{b_j})_n, & 0 < n \end{cases}$$

where

$$(\psi_{b_j})_n \stackrel{\text{def}}{=} \begin{cases} \begin{cases} 0, & 0 < n < m \\ (h_{b_2})_n, & n \geq m \end{cases}, & j = 1 \\ \begin{cases} 2\bar{\lambda} (\bar{a}_3 h_{b_3}^{(\infty)} + \bar{b}_3 h_{a_3}^{(\infty)})_n, & 0 < n < m \\ (2\bar{\lambda} (\bar{a}_3 h_{b_3} + \bar{b}_3 h_{a_3}) + 2\bar{a}_3 \bar{b}_3 h_\lambda)_n, & n \geq m \end{cases}, & j = 2 \\ \begin{cases} - \left(2\bar{a}_3 \bar{b}_3 h_{a_2}^{(\infty)} + 2(\bar{a}_3 \bar{b}_3 + \bar{a}_2 \bar{b}_3) h_{a_3}^{(\infty)} + (\bar{a}_3^2 + 2\bar{a}_2 \bar{a}_3) h_{b_3}^{(\infty)} \right)_n, & 0 < n < m \\ - \left(2\bar{a}_3 \bar{b}_3 h_{a_2} + 2(\bar{a}_3 \bar{b}_2 + \bar{a}_2 \bar{b}_3) h_{a_3} + \bar{a}_3^2 h_{b_2} + 2\bar{a}_2 \bar{a}_3 h_{b_3} \right)_n, & n \geq m \end{cases}, & j = 3. \end{cases}$$

Using Corollary 4.4, for each $n = 0, \dots, m$, we can easily compute upper bounds $\hat{\psi}_{a_j}, \hat{\psi}_{b_j} \geq 0$ such that

$$|(\psi_{a_j})_n| \leq (\hat{\psi}_{a_j})_n \quad \text{and} \quad |(\psi_{b_j})_n| \leq (\hat{\psi}_{b_j})_n, \quad \text{for all } n = 0, \dots, m.$$

More explicitly, we set $(\hat{\psi}_1^a)_n = (\hat{\psi}_1^b)_n = 0$ for $n = 0, \dots, m$, and

$$\begin{aligned} (\hat{\psi}_{a_2})_n &= 2|\bar{\lambda}| \Psi_n(\bar{a}_3), \\ (\hat{\psi}_{a_3})_n &= 2\Psi_n(\bar{a}_2 \bar{a}_3) + \Psi_n(\bar{a}_3^2), \\ (\hat{\psi}_{b_2})_n &= 2|\bar{\lambda}| (\Psi_n(\bar{a}_3) + \Psi_n(\bar{b}_3)), \\ (\hat{\psi}_{b_3})_n &= 2\Psi_n(\bar{a}_3 \bar{b}_3) + 2\Psi_n(\bar{a}_3 \bar{b}_2 + \bar{a}_2 \bar{b}_3) + \Psi_n(\bar{a}_3^2) + 2\Psi_n(\bar{a}_2 \bar{a}_3). \end{aligned}$$

Using these bounds, for $n = 0, \dots, m-1$, we compute $(\hat{z}_{a_j})_n$ and $(\hat{z}_{b_j})_n$ such that $|(z_{a_j})_n| \leq (\hat{z}_{a_j})_n$ and $|(z_{b_j})_n| \leq (\hat{z}_{b_j})_n$. More explicitly, for each $j = 1, 2, 3$, set

$$(\hat{z}_{a_j})_0 = (\hat{z}_{b_j})_0 = \frac{1}{\nu},$$

and

$$(\hat{z}_{a_j})_n \stackrel{\text{def}}{=} (|T|\hat{\psi}_{a_j})_n \quad (\hat{z}_{b_j})_n \stackrel{\text{def}}{=} (|T|\hat{\psi}_{b_j})_n,$$

for $n = 1, \dots, m-1$, where $|T|$ represents the operator with entries given by the component-wise absolute values of the entries of T . Moreover, set

$$\hat{z}_\lambda \stackrel{\text{def}}{=} 0, \quad \hat{z}_\delta = \frac{1}{\nu}, \quad \hat{z}_\gamma = \frac{1}{\nu}.$$

Recall that $z \stackrel{\text{def}}{=} [DF(\bar{x}) - A^\dagger]h$. Denote $w = Az$ and

$$w = (w_\lambda, w_\delta, w_{a_1}, w_{a_2}, w_{a_3}, w_\gamma, w_{b_1}, w_{b_2}, w_{b_3}).$$

Thus, for each $\alpha \in \{a_1, a_2, a_3, b_1, b_2, b_3\}$, we have the estimate

$$\begin{aligned} \|w_\alpha\|_\nu &= \sum_{n=0}^{m-1} \left| \left((A^{(m)} z^{(m)})_\alpha \right)_n \right| \omega_n + \sum_{n \geq m} \frac{1}{2n} |(w_\alpha)_n| \omega_n \\ &\leq \sum_{n=0}^{m-1} \left((|A^{(m)}| \hat{z}^{(m)})_\alpha \right)_n \omega_n + \frac{1}{2m} \sum_{n \geq m} |(T\psi_\alpha)_n| \omega_n \\ &\leq z_\alpha^{(1)} \stackrel{\text{def}}{=} \sum_{n=0}^{m-1} \left((|A^{(m)}| \hat{z}^{(m)})_\alpha \right)_n \omega_n + \frac{1}{2m} \|T\|_{B(\ell_\nu^1)} \hat{\psi}_\alpha^{(\infty)}, \end{aligned}$$

where

$$\hat{\psi}_\alpha^{(\infty)} \stackrel{\text{def}}{=} \begin{cases} 2\bar{\lambda}\|\bar{a}_3\|_\nu + \|\bar{a}_3^2\|_\nu, & \alpha = a_2 \\ 2\|\bar{a}_2\bar{a}_3\|_\nu + \|\bar{a}_3^2\|_\nu, & \alpha = a_3 \\ 2\bar{\lambda}(\|\bar{a}_3\|_\nu + \|\bar{b}_3\|_\nu) + 2\|\bar{a}_3\bar{b}_3\|_\nu, & \alpha = b_2 \\ 2\|\bar{a}_3\bar{b}_3\|_\nu + 2\|\bar{a}_3\bar{b}_2 + \bar{a}_2\bar{b}_3\|_\nu + \|\bar{a}_3^2\|_\nu + 2\|\bar{a}_2\bar{a}_3\|_\nu, & \alpha = b_3. \end{cases}$$

Therefore, we set

$$(4.26) \quad Z_1 = \max \left\{ z_\lambda^{(1)}, z_\delta^{(1)}, z_{a_1}^{(1)}, z_{a_2}^{(1)}, z_{a_3}^{(1)}, z_\gamma^{(1)}, z_{b_1}^{(1)}, z_{b_2}^{(1)}, z_{b_3}^{(1)} \right\},$$

where

$$z_\lambda^{(1)} \stackrel{\text{def}}{=} \left(|A^{(m)}| \hat{z}^{(m)} \right)_\lambda, \quad z_\delta^{(1)} \stackrel{\text{def}}{=} \left(|A^{(m)}| \hat{z}^{(m)} \right)_\delta, \quad z_\gamma^{(1)} \stackrel{\text{def}}{=} \left(|A^{(m)}| \hat{z}^{(m)} \right)_\gamma.$$

4.6. Z_2 bound. We look for a bound Z_2 in (4.15) such that

$$\|A[DF(\bar{x} + z) - DF(\bar{x})]\|_{B(X)} \leq Z_2 r, \quad \forall z \in B_r(0).$$

Let $z \in B_r(0)$ and $h \in B_1(0)$ (which we denote component-wise as in (4.25)), and denote

$$w = w(z, h) \stackrel{\text{def}}{=} (DF(\bar{x} + z) - DF(\bar{x}))h.$$

Note that $w_\lambda = w_\delta = w_\gamma = 0$, $w_{a_1} = w_{b_1} = 0 \in \ell_\nu^1$. Moreover, for each $\alpha \in \{a_2, a_3, b_2, b_3\}$, we have $w_\alpha = T\hat{w}_\alpha$ with

$$\begin{aligned} \hat{w}_{a_2} &= 2\bar{a}_3 h_\lambda z_{a_3} + 2h_{a_3} \bar{\lambda} z_{a_3} + h_\lambda z_{a_3}^2 + 2\bar{a}_3 h_{a_3} z_\lambda + 2h_{a_3} z_{a_3} z_\lambda, \\ \hat{w}_{a_3} &= -2\bar{a}_3 h_{a_3} z_{a_2} - 2\bar{a}_3 h_{a_2} z_{a_3} - 2\bar{a}_2 h_{a_3} z_{a_3} - 2h_{a_3} z_{a_2} z_{a_3} - h_{a_2} z_{a_3}^2, \\ \hat{w}_{b_2} &= 2\bar{b}_3 h_\lambda z_{a_3} + 2h_{b_3} \bar{\lambda} z_{a_3} + 2\bar{a}_3 h_\lambda z_{b_3} + 2h_{a_3} \bar{\lambda} z_{b_3} \\ &\quad + 2h_\lambda z_{a_3} z_{b_3} + 2\bar{b}_3 h_{a_3} z_\lambda + 2\bar{a}_3 h_{b_3} z_\lambda + 2h_{b_3} z_{a_3} z_\lambda + 2h_{a_3} z_{b_3} z_\lambda, \\ \hat{w}_{b_3} &= -2\bar{b}_3 h_{a_3} z_{a_2} - 2\bar{a}_3 h_{b_3} z_{a_2} - 2\bar{b}_3 h_{a_2} z_{a_3} - 2\bar{b}_2 h_{a_3} z_{a_3} - 2\bar{a}_3 h_{b_2} z_{a_3} \\ &\quad - 2\bar{a}_2 h_{b_3} z_{a_3} - 2h_{b_3} z_{a_2} z_{a_3} - h_{b_2} z_{a_3}^2 - 2\bar{a}_3 h_{a_3} z_{b_2} - 2h_{a_3} z_{a_3} z_{b_2} \\ &\quad - 2\bar{a}_3 h_{a_2} z_{b_3} - 2\bar{a}_2 h_{a_3} z_{b_3} - 2h_{a_3} z_{a_2} z_{b_3} - 2h_{a_2} z_{a_3} z_{b_3}. \end{aligned}$$

Thus

$$\begin{aligned} \|\hat{w}_{a_2}\|_\nu &\leq (4\|\bar{a}_3\|_\nu + 2|\bar{\lambda}| + 3r) r, \\ \|\hat{w}_{a_3}\|_\nu &\leq (4\|\bar{a}_3\|_\nu + 2\|\bar{a}_2\|_\nu + 3r) r, \\ \|\hat{w}_{b_2}\|_\nu &\leq (4\|\bar{b}_3\|_\nu + 4\|\bar{a}_3\|_\nu + 4|\bar{\lambda}| + 6r) r, \\ \|\hat{w}_{b_3}\|_\nu &\leq (4\|\bar{a}_2\|_\nu + 2\|\bar{b}_2\|_\nu + 4\|\bar{b}_3\|_\nu + 8\|\bar{a}_3\|_\nu + 9r) r. \end{aligned}$$

Note that $\|T\|_{B(\ell_\nu^1)} \leq 2\nu$, because

$$\begin{aligned} \|Th\|_\nu &= 2 \sum_{j \geq 1} | -h_{j-1} + h_{j+1} | \nu^j \leq \nu \left(2 \sum_{j \geq 1} |h_{j-1}| \nu^{j-1} \right) + \frac{1}{\nu} \left(2 \sum_{j \geq 1} |h_{j+1}| \nu^{j+1} \right) \\ &= \nu (|h_0| + \|h\|_\nu) + \frac{1}{\nu} (\|h\|_\nu - 2|h_1| - |h_0|) \\ &= \left(\nu + \frac{1}{\nu} \right) \|h\|_\nu + \left(\nu - \frac{1}{\nu} \right) |h_0| - \frac{2}{\nu} |h_1| \\ &\leq \left(\nu + \frac{1}{\nu} \right) \|h\|_\nu + \left(\nu - \frac{1}{\nu} \right) \|h\|_\nu = (2\nu) \|h\|_\nu. \end{aligned}$$

Fix $r^* \geq r$ (a condition that needs to be checked a posteriori), and set

$$\begin{aligned} z_{a_2}^{(2)} &\stackrel{\text{def}}{=} 4\|\bar{a}_3\|_\nu + 2|\bar{\lambda}| + 3r^*, \\ z_{a_3}^{(2)} &\stackrel{\text{def}}{=} 4\|\bar{a}_3\|_\nu + 2\|\bar{a}_2\|_\nu + 3r^*, \\ z_{b_2}^{(2)} &\stackrel{\text{def}}{=} 4\|\bar{b}_3\|_\nu + 4\|\bar{a}_3\|_\nu + 4|\bar{\lambda}| + 6r^*, \\ z_{b_3}^{(2)} &\stackrel{\text{def}}{=} 4\|\bar{a}_2\|_\nu + 2\|\bar{b}_2\|_\nu + 4\|\bar{b}_3\|_\nu + 8\|\bar{a}_3\|_\nu + 9r^*. \end{aligned}$$

Under these assumptions, we can verify that

$$\|A[DF(\bar{x} + z) - DF(\bar{x})]\|_{B(X)} \leq \|A\|_{B(X)} \|T\|_{B(\ell_\nu^1)} \max_{\alpha \in \{a_2, a_3, b_2, b_3\}} \{\|\hat{w}_\alpha\|_\nu\}.$$

Therefore, we can set

$$(4.27) \quad Z_2 \stackrel{\text{def}}{=} 2\nu \|A\|_{B(X)} \max \left(z_{a_2}^{(2)}, z_{a_3}^{(2)}, z_{b_2}^{(2)}, z_{b_3}^{(2)} \right),$$

where the computation of $\|A\|_{B(X)}$ is obtained with the same approach in Section 4.4.

4.7. Proof of Theorem 1.2. We fix $m = 65$ and obtain (using Newton's method) a numerical approximation $\bar{x} \in \mathbb{R}^{6m+3} = \mathbb{R}^{393}$ such that $F^{(m)}(\bar{x}) \approx 0$. We fixed $\nu = 1.05$, and combining the explicit and computable bounds Y_0 , Z_0 , Z_1 and Z_2 given respectively by (4.21), (4.24), (4.26) and (4.27), we defined the radii polynomial $p(r)$ as in (4.16) and applied the radii polynomial approach of Theorem 4.1 to show that $p(r_0) < 0$ with $r_0 = 5.7 \times 10^{-12}$. This yields the existence of a unique $\tilde{x} \in B_{r_0}(\bar{x}) \subset X$ such that $F(\tilde{x}) = 0$. This rigorous error bound implies the proof of Theorem 1.2. The graph of the solution is portrayed in Figure 1.

Choosing the values for m and ν is heuristic, non unique and done essentially so that the bound Z_1 satisfies $Z_1 < 1$ (this is indeed a necessary condition for (4.17) to hold for some $r_0 > 0$). Recalling (4.26), and the definition $\hat{z}_\delta = \hat{z}_\gamma = \frac{1}{\nu}$, it is clear that $Z_1 < 1$ only if $\nu > 1$. However, it cannot be taken too large as the solution itself may not have enough regularity to be in the space ℓ_ν^1 . Also, taking $\nu > 1$ provides decay in the bound $\Psi_k(\bar{\alpha})$ in (4.23). Finally the choice of m needs to be large enough so that the *defect bound* Y_0 is small enough and so that the *tail terms* $\frac{1}{2m} \|T\|_{B(\ell_\nu^1)} \hat{\psi}_\alpha^{(\infty)}$ in the definition of $z_\alpha^{(1)}$ are less than 1. While other choices would have worked, we found that the choice $m = 65$ and $\nu = 1.05$ yielded the best rigorous error bound possible in that of $r_0 = 5.7 \times 10^{-12}$.

5. Conclusion. In this paper, we studied a particular one-dimensional model for a microelectrical device in equation (1.1). This equation forms part of a wide range of Hamiltonian PDEs modeling physical phenomena such as the nonlinear wave equation, the nonlinear Schrödinger equation, beam's equation, and Euler's equation and its multiple approximations appearing in water waves. All these Hamiltonian PDEs exhibit trivial or steady solutions, and near these steady solutions there are periodic and quasiperiodic solutions. However, proving existence of such solutions exhibits a small divisor problem unless one imposes a special relation between the period and the domain of the equation. In that case the equation can be solved with a Lyapunov-Schmidt procedure by separating the equation into the kernel and range equation. Even if the small divisor problem can be avoided, there are other mathematical difficulties associated to the existence of periodic solutions such as the lack of compactness of the linearized operator for the range equation or the infinite dimension of the kernel equation.

Equation (1.1) has a family of stable steady states (u_λ, λ) for $\lambda \in [0, \lambda_*]$. In this paper, we proved the existence of an infinite number of continuous branches of periodic solutions arising from the steady solution u_λ . The local branches have fixed periods satisfying a rational relation with the space length and arise from u_λ for bifurcation values $\lambda \in [0, \lambda_*]$. In order to tackle the difficulties associated to prove this fact we used a combination of analytic estimates and computer-assisted proofs. We also introduced a systematic setting to compute numerically these branches of solutions.

The specific features of equation (1.1) imply that the set of bifurcation values λ is not only infinite, but actually it is a dense subset of parameters in $[0, \lambda_*]$. Moreover, in the complement of the dense set of $[0, \lambda_*]$, KAM theory for Hamiltonian PDEs may be used to prove the existence of periodic solutions for a subset of almost full measure of $[0, \lambda_*]$. However, in such case the periods must satisfy some Diophantine relations and the periodic solutions form cantor-like sets, that is they do not form continuous families as in our results.

We finish our conclusion by mentioning some of the possible extensions of our work:

1. The methods presented here can be used to solve similar problems in other Hamiltonian PDEs. In particular, the same procedure can be implemented to prove the existence of periodic solutions in the equation of MEMS with other dielectric permittivity properties $f(x) \neq 1$ or for other nonlinear wave equations of the form $u_{tt} - u_{xx} = f(x, u)$.
2. In order to prove our result, in Section 4 we implement a computer-assisted proof to validate the steady solution u_{λ_*} at critical value λ_* . Actually, the numerical setting in Sections 3.1 and 3.2, and the methods of Section 4 can be used to validate also the trivial branch u_λ and its spectrum (eigenfunctions and eigenvalues). Furthermore, similar procedures can be used to validate also radial steady solutions for the nonlinear equation in more dimensions

$$\Delta u = \lambda f(u), \quad u_{\partial B} = 0,$$

with analytic nonlinearities such as $f(u) = u^p$ or $f(u) = e^u$ (e.g. see [3]).

3. Unfortunately, there are no readily available methods to validate numerically the periodic solutions obtained in Section 3.3. The problem is that the inverse of the hyperbolic operator L is only bounded but not compact. Indeed, the methods of Section 4 depend strongly on the compactness of the inverse operators to validate the numerical solutions, because the compactness allows to obtain estimates for the Galerkin approximation or truncation of the linear operators. Further research and new ideas are required to validate rigorously the numerical computations of the periodic solutions.

6. Acknowledgments. CGA is indebted to G. Flores and M. Tejada-Wriedt for discussions related to this project.

REFERENCES

- [1] H. Amann and E. Zehnder. Nontrivial solutions for a class of nonresonance problems and applications to nonlinear differential equations. *Ann. Scuola Norm. Sup. Pisa Cl. Sci. (4)*, 7(4):539–603, 1980.
- [2] G. Arioli and H. Koch. Families of periodic solutions for some Hamiltonian PDEs. *SIAM J. Appl. Dyn. Syst.*, 16(1):1–15, 2017.
- [3] I. Balázs, J. B. van den Berg, J. Courtois, J. Dudás, J.-P. Lessard, A. Vörös-Kiss, J. F. Williams,

- and X. Y. Yin. Computer-assisted proofs for radially symmetric solutions of PDEs. *J. Comput. Dyn.*, 5(1-2):61–80, 2018.
- [4] R. Castelli, M. Gameiro, and J.-P. Lessard. Rigorous numerics for ill-posed PDEs: periodic orbits in the Boussinesq equation. *Arch. Ration. Mech. Anal.*, 228(1):129–157, 2018.
- [5] P. H. Chang and H. A. Levine. The quenching of solutions of semilinear hyperbolic equations. *SIAM J. Math. Anal.*, 12(6):893–903, 1981.
- [6] W. Craig and C. García-Azpeitia. Standing waves of fixed period for $n + 1$ vortex filaments. *Journal of Dynamics and Differential Equations*, 2019.
- [7] W. Craig and C. E. Wayne. Newton’s method and periodic solutions of nonlinear wave equations. *Comm. Pure Appl. Math.*, 46(11):1409–1498, 1993.
- [8] M. G. Crandall and P. H. Rabinowitz. Bifurcation from simple eigenvalues. *Journal of Functional Analysis*, 8(2):321–340, oct 1971.
- [9] S. Day, J.-P. Lessard, and K. Mischaikow. Validated continuation for equilibria of PDEs. *SIAM J. Numer. Anal.*, 45(4):1398–1424 (electronic), 2007.
- [10] R. de la Llave. Variational methods for quasi-periodic solutions of partial differential equations. In *Hamiltonian systems and celestial mechanics (Pátzcuaro, 1998)*, volume 6 of *World Sci. Monogr. Ser. Math.*, pages 214–228. World Sci. Publ., River Edge, NJ, 2000.
- [11] G. Flores. Dynamics of a damped wave equation arising from MEMS. *SIAM J. Appl. Math.*, 74(4):1025–1035, 2014.
- [12] G. Flores, G. Mercado, and J. A. Pelesko. Dynamics and touchdown in electrostatic MEMS. In *Proceedings International Conference on MEMS, NANO and Smart Systems, Banff, Alberta, Canada*, pages 182–187, 2003.
- [13] G. Flores, G. Mercado, J. A. Pelesko, and N. Smyth. Analysis of the dynamics and touchdown in a model of electrostatic MEMS. *SIAM J. Appl. Math.*, 67(2):434–446, 2006/07.
- [14] M. Gameiro and J.-P. Lessard. Rigorous computation of smooth branches of equilibria for the three dimensional Cahn-Hilliard equation. *Numer. Math.*, 117(4):753–778, 2011.
- [15] C. García-Azpeitia. Standing waves in a counter-rotating vortex filament pair. *J. Differential Equations*, 264(6):3918–3932, 2018.
- [16] N. Ghoussoub and Y. Guo. On the partial differential equations of electrostatic MEMS devices: stationary case. *SIAM J. Math. Anal.*, 38(5):1423–1449, 2006/07.
- [17] C. Gugg, T. J. Healey, H. Kielhöfer, and S. Maier-Paape. Nonlinear standing and rotating waves on the sphere. *J. Differential Equations*, 166(2):402–442, 2000.
- [18] Y. Guo. Dynamical solutions of singular wave equations modeling electrostatic MEMS. *SIAM J. Appl. Dyn. Syst.*, 9(4):1135–1163, 2010.
- [19] Y. Guo, Z. Pan, and M. J. Ward. Touchdown and pull-in voltage behavior of a MEMS device with varying dielectric properties. *SIAM J. Appl. Math.*, 66(1):309–338, 2005.
- [20] A. Hungria, J.-P. Lessard, and J. D. Mireles James. Rigorous numerics for analytic solutions of differential equations: the radii polynomial approach. *Math. Comp.*, 85(299):1427–1459, 2016.
- [21] N. I. Kavallaris, A. A. Lacey, C. V. Nikolopoulos, and D. E. Tzanetis. A hyperbolic non-local problem modelling MEMS technology. *Rocky Mountain J. Math.*, 41(2):505–534, 2011.
- [22] N. I. Kavallaris, A. A. Lacey, C. V. Nikolopoulos, and D. E. Tzanetis. On the quenching behaviour of a semilinear wave equation modelling MEMS technology. *Discrete Contin. Dyn. Syst.*, 35(3):1009–1037, 2015.
- [23] H. B. Keller. *Lectures on numerical methods in bifurcation problems*, volume 79 of *Tata Institute of Fundamental Research Lectures on Mathematics and Physics*. Published for the Tata Institute of Fundamental Research, Bombay, 1987. With notes by A. K. Nandakumaran and Mythily Ramaswamy.
- [24] H. Kielhöfer. Bifurcation of periodic solutions for a semilinear wave equation. *J. Math. Anal. Appl.*, 68(2):408–420, 1979.
- [25] H. Kielhöfer. A bifurcation theorem for potential operators. *J. Funct. Anal.*, 77(1):1–8, 1988.
- [26] J. M. Ortega. The Newton-Kantorovich theorem. *Amer. Math. Monthly*, 75:658–660, 1968.
- [27] J. A. Pelesko. Mathematical modeling of electrostatic MEMS with tailored dielectric properties. *SIAM Journal on Applied Mathematics*, 62(3):888–908, jan 2002.
- [28] P. H. Rabinowitz. Free vibrations for a semilinear wave equation. *Comm. Pure Appl. Math.*, 31(1):31–68, 1978.
- [29] K. Rektorys. *Variational Methods in Mathematics, Science and Engineering*. Springer Netherlands, 1977.
- [30] S. Rump. INTLAB - INTerval LABoratory. In T. Csendes, editor, *Developments in Reliable Computing*, pages 77–104. Kluwer Academic Publishers, Dordrecht, 1999. <http://www.ti3.tu-harburg.de/rump/>.
- [31] R. A. Smith. On a hyperbolic quenching problem in several dimensions. *SIAM J. Math. Anal.*,

- 20(5):1081–1094, 1989.
- [32] J. B. van den Berg and J.-P. Lessard. Rigorous numerics in dynamics. *Notices of the American Mathematical Society*, 62(9):1057–1061, 2015.
- [33] J. B. van den Berg, J.-P. Lessard, and K. Mischaikow. Global smooth solution curves using rigorous branch following. *Math. Comp.*, 79(271):1565–1584, 2010.